## Texts and Monographs in Physics

Series Editors:<br>R. Balian, Gif-sur-Yvette, France<br>W. Beiglböck, Heidelberg, Germany<br>H. Grosse, Wien, Austria<br>W. Thirring, Wien, Austria

Vladimir M. Akulin

# Coherent Dynamics of Complex Quantum Systems 

With 140 Figures

6) Springer

Professor Vladimir M. Akulin<br>Laboratoire Aimé Cotton<br>Bat 505<br>Campus d'Orsay<br>91405 Orsay Cedex, France

ISSN 0172-5998

## ISBN-10 3-540-21052-0 Springer Berlin Heidelberg New York ISBN-13 978-3-540-21052-8 Springer Berlin Heidelberg New York

This work is subject to copyright. All rights are reserved, whether the whole or part of the material is concerned, specifically the rights of translation, reprinting, reuse of illustrations, recitation, broadcasting, reproduction on microfilm or in any other way, and storage in data banks. Duplication of this publication or parts thereof is permitted only under the provisions of the German Copyright Law of September 9, 1965, in its current version, and permission for use must always be obtained from Springer. Violations are liable for prosecution under the German Copyright Law.
Springer is a part of Springer Science+Business Media
springeronline.com
(c) Springer-Verlag Berlin Heidelberg 2006

Printed in the Netherlands
The use of general descriptive names, registered names, trademarks, etc. in this publication does not imply, even in the absence of a specific statement, that such names are exempt from the relevant probreak tective laws and regulations and therefore free for general use.

Typesetting: Data conversion by LE-TeX Jelonek, Schmidt \& Vöckler GbR
Cover design: design \& production GmbH , Heidelberg
Printed on acid-free paper $\quad$ SPIN $10971628 \quad 55 / 3141 / \mathrm{mh} \quad 543210$

## Preface

This book is an attempt to put together a large number of similar problems that one encounters in different fields of modern quantum physics and that have common features considering multilevel quantum systems. The main motivation was to present from the same standpoints various models and approaches that have been developed in atomic, molecular, condensed matter, chemical, laser and nuclear physics in various contexts.

The book is based on my lectures in the Moscow Institute of Physics and Technology, in the Aime Cotton Laboratory of CNRS, and some other courses that I have delivered during last two decades. It includes the original results obtained in collaboration with my collegues V. Aquilanti, I. Averbukh, A. Belousov, M. Blaauboer, E. Borsela, E. Brion, B. Brunetti, C. Brechignac, P. Cahuzac, F. Carlier, I. Dumer, V. Gershkovich, G. Esadze, G. Garsevanishvili, G. Harel, E. Khokhlov, G. Kurizki, R. Larciprete, I. Murachko, A. Nesterenko, A. Orlov, S. Pelegrin, P. Pillet, F. Rebentrost, A. Sarfati, E. Schlag, W. Schleich, F. Vecchiocativi from different scientific centres in the world and with whom I had the pleasure to work on the dynamical aspects of the behaviour of complex quantum systems. I express my deep gratitude to them for their collaboration. The book also contains numerous results of other authors that have, however, been expounded in different notations consistent with the present text, and that sometimes even rely on an alternative derivation as compared to the original version.

In preparing the text I decided to add several results both scientific, yet unpublished, and pedagogical that I feel are necessary for giving the entire picture of the processes in complex quantum systems. Some of these results, presented in Chap. 6, have been obtained in collaboration with V. Kravtsov to whom I express my sincere acknowledgments. I also very much appreciate the discussions with V. Agranovich, E. Bogomolny, B. Chirikov, T. Gallagher, P. Golovinskii, M. Fedorov, Y. Fyodorov, C. Jungen, J. Jortner, L. Maksimov, V. Man’ko, I. Mazets, V. Kac, A. Kazantsev, D. Khmelnitski, A. Kofman, I. Lerner, E. Nikitin, M. Shapiro, D. Shepelyanski, V. Pokrovskii and A. Prokhorov on the different aspects of the results presented in the book.

I want to worship the memory of G. Askaryan who has influenced my choice of profession, showing me the beauty of physics presented in our life.

Finally, I express my profound gratitude to my teachers Alexandr Dyknne and Nikolay Karlov.

Paris, May 2004
Vladimir Akulin

## Contents

1 Complex Systems and Their Statistical Description ..... 1
2 Examples of Complex Systems ..... 17
2.1 Molecules and Atoms in Laser Fields ..... 18
2.1.1 Laser Breaking of a Weakly Bonded Complex ..... 18
2.1.2 Laser-Induced Electronic Transitions in Molecules ..... 21
2.1.3 Vibrational Excitation of Polyatomic Molecules ..... 22
2.1.4 Transitions Among Levels with Fine Structure ..... 24
2.1.5 Excitation of Rydberg States in Atoms ..... 26
2.1.6 Competition of Multiphoton Processes of Different Orders ..... 27
2.2 Collisions and Reactions of Molecules ..... 28
2.2.1 Collisional Redistribution of Energy ..... 28
2.2.2 Chemical Reactions ..... 32
2.2.3 Intermolecular Conversion and Photochemistry ..... 35
2.3 Rydberg Molecules ..... 38
2.3.1 Subthreshold Photoionization ..... 38
2.3.2 Collisional Ionization ..... 40
2.4 Atomic and Molecular Clusters ..... 41
2.4.1 Ground Electronic State of Hot Metallic Clusters ..... 43
2.4.2 Optical Properties of Clusters ..... 46
2.5 Some Other Examples ..... 48
2.5.1 Ion Traps ..... 48
2.5.2 Disordered Solids and Surfaces ..... 51
2.5.3 Nonlinear Optics ..... 55
2.5.4 Cooperative Effect ..... 56
2.5.5 Many-Body Effects in Cold Rydberg Gas ..... 58
3 Two-Level and Level-Band Systems ..... 61
3.1 Two-Level System ..... 62
3.2 Level-band System ..... 68
3.2.1 General Consideration ..... 68
3.2.2 Continuous Band Model ..... 70
3.2.3 Measurements and Relaxation as Processes in Level--Continuum Systems ..... 83
3.3 Long-Time Behavior ..... 88
3.3.1 General Consideration of the Long-Time Limit ..... 88
3.3.2 Quantum Recurrences ..... 91
3.3.3 Quantum Revivals ..... 95
3.3.4 Fractional Revivals ..... 99
3.3.5 Revivals and the Classical Limit ..... 104
3.4 Population of Inhomogeneous Bands ..... 105
3.4.1 Statistically Independent Levels ..... 106
3.4.2 Factorization of the Level Population and the Ensemble Average ..... 108
3.4.3 The Long-Time Asymptotic ..... 111
4 Two-Band System ..... 123
4.1 General Consideration ..... 124
4.1.1 Series and Diagrams for the Level-Band Problem ..... 124
4.1.2 Series and Diagrams for the Two-Band Problem ..... 128
4.1.3 The Renormalization ..... 133
4.2 Non-Degenerate Bands ..... 139
4.2.1 General Remarks and the Main Questions ..... 139
4.2.2 Renormalized Energies and the Population Distribution ..... 140
4.2.3 Dynamics of the Total Populations of Bands ..... 143
4.2.4 Different Bands ..... 144
4.3 Two Degenerate Levels ..... 145
4.3.1 Degenerate Levels as a Complex System ..... 145
4.3.2 The Bands as an Ensemble of Two-Level Systems ..... 148
4.4 A Band Coupled to a Degenerate Level ..... 151
4.4.1 Total Population of the Bands ..... 151
4.4.2 Population Distribution over the Band ..... 154
4.4.3 Role of the Interaction Rank ..... 159
4.5 The Role of Correlations ..... 161
4.5.1 Two levels and a band ..... 161
4.5.2 Two Bands With a Correlated Coupling ..... 169
4.5.3 Regime of Stabilization for the Correlated Coupling ..... 178
4.5.4 Correlation Between the Mean Squared Coupling and the Energy Position ..... 179
5 Soluble Time-Dependent Systems ..... 187
5.1 Algebraic Structure of Time Dependent Systems ..... 188
5.2 Time-Dependent Two-Level Systems ..... 194
5.2.1 Landau-Zener Problem ..... 194
5.2.2 Landau-Zener Transition to a Decaying State ..... 199
5.2.3 Landau-Zener Transition in the Presence of Transversal Relaxation ..... 202
5.2.4 Excitation by a Pulse ..... 208
5.2.5 Exponentially Rising Coupling ..... 210
5.3 Semiclassical Analysis of Time-Dependent Systems ..... 213
5.3.1 Two-Level Systems and the Dykhne Formula ..... 213
5.3.2 Multilevel Systems ..... 218
5.4 Time-Dependent Level-Band System ..... 221
5.4.1 The Demkov-Osherov Problem ..... 222
5.4.2 The Landau-Zener Transition at the Continuum Edge ..... 227
6 Time-Dependent Complex Systems ..... 237
6.1 Degenerate Level Crosses an Infinite Band ..... 238
6.2 Perturbation Proportional to a Random Matrix ..... 244
6.2.1 Population Distribution ..... 246
6.2.2 Response to Perturbation Proportional to a Random Matrix ..... 250
6.3 Harmonic Perturbation of Complex Systems ..... 261
6.3.1 Population Distribution over a Uniform Spectrum ..... 261
6.3.2 Response of the Uniform Spectrum to a Harmonic Perturbation ..... 265
6.4 Two-Frequency Excitation of Complex Systems ..... 270
6.4.1 Population Dynamics for Bi-Harmonic Excitation ..... 270
6.4.2 Response to Bi-Harmonic Excitation ..... 273
6.5 Two-Band System in a Periodic Field ..... 277
6.5.1 Dynamics of Total Band Populations ..... 277
6.5.2 Population Distribution over the Bands ..... 282
6.6 Control of Complex Quantum Systems ..... 284
6.6.1 Control of Two-Level Systems ..... 286
6.6.2 Holonom and Non-Holonom Systems ..... 294
6.6.3 Control of Coherence Loss ..... 302
7 The Dynamics of One-Dimensional Relay-Type Systems ..... 309
7.1 Exactly Soluble Relays of Isolated Levels ..... 309
7.1.1 Uniform Coupling and Linear Detuning ..... 310
7.1.2 The Harmonic Oscillator in an Arbitrary Time-Dependent Field ..... 311
7.1.3 Raman Pumping of a Harmonic Oscillator ..... 313
7.1.4 The Harmonic Oscillator in the Simultaneous Presence of Dipole and Raman Pumping ..... 316
7.2 General Case of an Exactly Soluble Relay ..... 317
7.2.1 Conditions for the Existence of a Polynomial Solution ..... 318
7.2.2 The Increasing Coupling $\left|V_{n}\right|=\sqrt{a(n-b)(n-c)}$. ..... 322
7.2.3 Decreasing Coupling $\left|V_{n}\right|=1 / \sqrt{a n+b}$ ..... 324
7.3 Smooth Variation of the Parameters ..... 325
7.3.1 WKB approximation ..... 326
7.3.2 Position and Width of the Erenfest Wavepacket ..... 327
7.3.3 The Tunneling Probability ..... 329
7.3.4 Applicability of the WKB Approximation ..... 330
7.4 Relay with disordered parameters ..... 331
7.4.1 Ensemble Averaged Amplitudes and Corresponding Populations ..... 333
7.4.2 Ensemble Averaged Spectrum ..... 335
7.4.3 Distribution of the Amplitude Ratios ..... 335
7.4.4 Distribution of the Populations for Long Times ..... 340
7.4.5 Dynamics of the Asymptotic Populations ..... 342
7.5 Field Theory Method for Disordered Systems ..... 344
7.5.1 Tunneling Transparency and Classical Bosonic Fields ..... 344
7.5.2 An Analogy With the Liouville Equation ..... 352
7.5.3 Classical Fermionic Fields for the Population Dynamics ..... 354
7.6 Population Dynamics in a Disordered Chain ..... 360
7.6.1 Population Dynamics and Propagating Fictitious Particles ..... 360
7.6.2 Mapping over a Period and the Ensemble Average ..... 364
7.6.3 Time Dependence of the Population Distribution ..... 369
8 Composite Complex Quantum Systems ..... 373
8.1 Relay of Multilevel Bands ..... 373
8.1.1 Degenerate Bands With Random Coupling ..... 376
8.1.2 Non-Degenerate Bands With Random Coupling ..... 380
8.1.3 Correlated Coupling ..... 387
8.2 Random Walks and Coherent Behavior ..... 392
8.2.1 Level Decay to a Band of Random Walks ..... 392
8.2.2 Interference of Random Returns at Long Times. General Consideration ..... 397
8.2.3 Three Types of Random Walk and the Asymptotic Decay ..... 406
8.3 Manifestation of Quantum Complexity in the State Density ..... 413
8.3.1 Spectrum Transformation Induced by Random Perturbation ..... 414
8.3.2 The Effect of Quantum Recurrences on the State Density Profiles ..... 422
8.3.3 The Density of Quantum States of Fractals ..... 431
9 Bibliography and Problems ..... 439
9.1 Bookshelf ..... 440
9.2 Problems ..... 449
References ..... 459
Index ..... 467

## 1 Complex Systems and Their Statistical Description

We often meet in physics various quantum systems comprising many quantum levels. One may say that the majority of problems of contemporary interest are of this type. Statistics is a natural way to describe large systems. Initially, the statistical approach to individual complex quantum objects had been formulated in the context of the interpretation of nuclear spectra where the average positions of the energy levels and their average relative distances were the principal concern. Later, the effects of complexity related to disorder were studied in solid state physics and this gave rise to the theory of dynamical localization, focused on the structure of wavefunctions. The rapid development of laser physics, the molecular beam technique, experiments with cold atoms, femtosecond physics, quantum optics and quantum informatics has formulated a series of theoretical problems centered on the time evolution of multilevel quantum objects. Analytical description of the dynamics of multilevel quantum systems, including the results which can be obtained with the help of the statistical approach, is the main subject of this book.

The dynamics of multilevel systems is a complicated process, which involves population exchange among many quantum states. Though a few types of multilevel problems can be solved analytically, the Schrödinger equation for a multilevel quantum system typically does not posses an exact analytical solution in terms of elementary or special functions nor in terms of integrals over a reasonably small number of variables. Therefore the exact solution implies the numerical integration of the Schrödinger equation. However, the numerical solution of such problems usually encounters difficulties of three different kinds. The first one is rather evident: as the number of states in a system increases, the numerical approach becomes awkward, since it becomes more and more difficult to handle the increasing amount of information. In some cases, this difficulty may arise even earlier, when in order to simply find the matrix elements of the Hamiltonian and to write down the Schrödinger equation in a finite basis representation, one has to perform complicated multidimensional integrations. The second difficulty arises due to the influence of higher-order resonances which may manifest themselves in a hierarchy of interactions that are different by many orders of magnitude. This makes the dynamics of a multilevel quantum system considerably different at different
time-scales. When the influence of the processes developing at exponentially different time-scales becomes important, the numerical approach fails to select an appropriate step-size for the computational algorithms.

There is an inevitable difficulty of the third kind. Provided one has obtained a numerical solution for a large multilevel quantum system and now knows all of the values of the quantum amplitudes for each state and for each moment of time, how is one to gain insight into this mass of information? How does one reach a consistent conclusion, address the main task of theoretical description and find a correspondence law between the input and the output of the calculations? Evidently, the physically meaningful quantities are just a few cumulative characteristics of the quantum system. But how can one identify these characteristics and find comprehensive relations among them? One may say that here we face a conflict between information and insight: the more information we have the more difficult it is to derive physically interesting conclusions.

## Complexity and the Ensemble Average

A natural question arises when an increasing number of states in a system renders the numerical treatment inefficient; do we really need to know all of the details of the dynamics of a multilevel quantum system as an output of the calculations, or should we rather look for its general properties by considering ensembles of many similar complex quantum objects? The general properties do not depend on all of the tiny details of the particular realization, but only on a small number of "crude" characteristics of quantum systems. We therefore arrive at the idea of the ensemble average, which after E.Wigner, became the key tool of analysis of complex quantum systems. This is in complete analogy with the thermodynamic approach, which summarizes general properties of macroscopically identical systems by ignoring their microscopic differences.

Historically, the random matrix approach to complex quantum systems and the philosophy of ensemble averages associated with this approach were first suggested by E.Wigner for the analysis of the spectra of complex nuclei. In fact, little is known about systems of many interacting nucleons except some general properties such as symmetry with respect to time inversion, permutations of the particles, and other transformations in phase space. Detailed information relating to the structure of the wavefunctions is hidden in the complexity of the system. Nature affords no opportunity to obtain such detailed information and thus provide insight into this structure in a reasonably simple way. The natural approach of mathematical formalization of our unavoidable ignorance is to assume that all matrix elements of the Hamiltonian $H_{i, j}$ are random numbers that satisfy just a small number of general symmetry conditions such as

$$
\begin{equation*}
H_{i, j}=H_{j, i} \quad \text { or } \quad H_{i, j}=H_{j, i}^{*}, \tag{1.1}
\end{equation*}
$$

for real or complex-valued Hamiltonians, respectively. In traditional calculations, the first step is the determination of the eigenvalues of the Hamiltonians. The logical generalization of this concept to the case of a complex system is to employ the distribution function of the eigenvalues. Three principal statistical distributions are usually considered in the analysis of the spectral properties of complex nuclei and their comparison to experimental data: (1) distribution of the eigenvalues, namely the mean density of states $g(E) \equiv \Delta_{1}(E) ;(2)$ distribution of the energy spacing between nearest neighbors $\Delta_{2}\left(E_{n}-E_{n-1}\right)$; and (3) the statistics of the deviations $\Delta_{3}\left(E_{n}-n \bar{g}\right)$ of the energy levels $E_{n}$ from the positions $\bar{g} n$ suggested by the mean local density of states $\bar{g}$. These universal characteristics are usually employed in other fields of physics whenever the quantum system in consideration possesses a complex spectrum.

The idea to employ the ensemble average for the description of complex quantum systems gave rise to a large branch of contemporary theoretical physics and related mathematics, that investigates eigenvalue statistics and the properties of the eigenfunctions of complex Hamiltonians. This approach has been successfully adopted in many areas of nuclear physics, in solid state physics, in spectroscopy, and the theory of quantum chaos. In atomic, molecular, and laser physics, a series of theoretical problems has also been formulated where the time evolution of a complex quantum object is the main interest. The frequency domains of the latter range from seconds for cold atoms to femtoseconds for intramolecular physics. The large number of quantum states participating in the dynamical processes allows one to use the ensemble average idea for the analysis of the behavior of these systems as well. However, the main physical characteristics to be found for these systems, such as the net population of a group of levels or the time dependence of optical polarization, differ essentially from the quantities that one calculates in nuclear or in solid state physics, where the emphasis is usually upon the statistics of eigenstates or on the typical spatial structure of eigenfunctions.

We should clearly understand the basic assumptions underlying the statistical approach, in order to avoid possible misunderstandings concerning the question of what we can and what we cannot pretend to describe with the help of this method. The essence of the ensemble average idea is the following: we are interested in properties that are not sensitive to the details or to the particular realization of a system, and therefore we consider ensembles of many systems that are different in the tiny details but have the same "crude" characteristics. We calculate the physical quantities of interest for each of these systems and take the average over the ensemble. In this way we find the common properties, while individual features of the systems are washed out. We note that technically there is no need to perform the average after the calculations of physical quantities of interest for each of the systems. On the contrary, it is much more convenient to do it at an earlier stage, before the summation of the all-order perturbation series, for example.

This usually leads to a considerable simplification of the calculations, that can therefore often be performed analytically.

The choice of ensemble is the most delicate matter of this approach. It mainly relies on our physical intuition. We have to decide a priori which details are important and which are not important for the "crude" properties of systems, and then we have to construct the corresponding ensemble. If, after the average, the mean value of the quantity of interest is not zero and its dispersion vanishes, one can say that this quantity is indeed "crude" and moreover the ensemble is chosen properly.

## Statistical Ensembles

In our consideration we employ a relatively general model of quantum systems, for which the Hamiltonian is not a completely random matrix but consists of two parts: one, $\widehat{H}_{0}$, is regular and the other, the perturbation $\widehat{V}$, is random. The matrix of the random part belongs to the so-called Gaussian unitary ensemble where both the real and the imaginary parts of the Hermitian matrix elements obey Gaussian statistics, that is they are normally distributed around the zero mean value. We shall not give here any proofs that this is the case for each particular problem, but just note that this natural assumption is in the spirit of the central limit theorem of statistics. We also note that this assumption usually allows us to obtain exact analytical expressions for many physical quantities of interest. Therefore this model may serve as a guide for qualitative understanding of the properties of other multilevel systems, for which the applicability of the statistical assumption is questionable, but the exact analytical solution is not available. Though later on we will be mainly concerned with the case where $\widehat{H}_{0}$ is regular and only $\widehat{V}$ is random, while all of the physical properties of the complex system result from the interplay of these two parts, here we briefly present the main classical results of random matrix theory that ignore the regular part and treat the Hamiltonian $\widehat{H}=\widehat{H}_{0}+\widehat{V}$ of a complex system as a completely random matrix.

Apart from the Gaussian unitary ensemble $H_{i, j}=H_{j, i}^{*}$, two other standard distributions with Gaussian statistics of matrix elements are often employed in the cases where a certain residual symmetry of the system considered is important. The Gaussian orthogonal ensemble $H_{i, j}=H_{j, i}$ is applicable to time-reversible systems. This implies that all of the matrix element are real, whereas the Gaussian simplectic ensemble $H_{1 i, 1 j}=H_{2 j, 2 i}^{*}$, $H_{1 i, 2 j}=-H_{2 j, 1 i}^{*}$ corresponds to the situation where each matrix consists of four blocks, with the diagonal blocks complex conjugated, and the offdiagonal blocks complex conjugated with inverted signs. The last ensemble describes the time-reversible situation in the absence of central symmetry.

All three Gaussian ensembles result in the same distribution of the eigenvalues of the random matrices with mean density $g(E)$ of the eigenstates given by the Wigner semicircle distribution

$$
\begin{equation*}
g(E)=\frac{\sqrt{E-N\left\langle V^{2}\right\rangle}}{\pi\left\langle V^{2}\right\rangle} \tag{1.2}
\end{equation*}
$$

where $N$ is the total number of states and $\left\langle V^{2}\right\rangle$ is the mean squared element of the random matrix. The main difference between these ensembles manifests itself in the statistics of the energy distances $E_{n+1}-E_{n}$ between the neighboring levels. The probability to find the nearest neighbor on the small energy distance scales as $E_{n+1}-E_{n},\left(E_{n+1}-E_{n}\right)^{2}$, and $\left(E_{n+1}-E_{n}\right)^{4}$ for the orthogonal, unitary, and simplectic ensembles, respectively. We will not consider here the physical properties associated with this difference but just note that all of these ensembles are "rigid" and small energy intervals between the neighboring levels are unlikely. In other words, the levels of complex systems interact and repel each other.

One more important case of the level statistics has to be mentioned. If a system is composed of two or more independent parts that do not interact among themselves, the spectrum of the compound is the direct superposition of the spectra of the parts. In this case, the neighboring levels may belong to different parts. They do not interact and do not repel each other. Therefore the probability to find levels at a small distance is equal to the mean density of states in this energy domain. Such a situation corresponds to the so-called Poisson distribution, typical of Markovian random processes.

We also note another approach to complex quantum systems which concentrates not on the statistics of the eigenvalues of random Hamiltonians, but on that of the evolution operators $\widehat{U}=\exp (-i t \widehat{H})$ governing the time evolution of the systems. In contrast to $\widehat{H}$, the eigenvalues of $\widehat{U}$ are not real but occupy the unit circle in the complex plane. Therefore the statistical ensembles introduced for these quantities by Dyson are called circular ensembles, and these ensembles can also be orthogonal, unitary, and simplectic by analogy to the Gaussian ensembles. Note that the circular ensembles had been originally introduced in nuclear physics, not for the time-dependent evolution operator, but for the scattering matrix which is roughly speaking the operator of evolution during a finite interaction time interval. The probability $P$ to find the eigenvalues of the scattering matrix at the positions $\left\{e^{-i \varphi_{n}}\right\}$ on the unit circle reads

$$
\begin{equation*}
P\left(\left\{e^{-i \varphi_{n}}\right\}\right) \sim \prod_{k<l}\left|\exp \left(-i \varphi_{k}\right)-\exp \left(-i \varphi_{l}\right)\right|^{\beta} \prod_{n=1}^{N} d \varphi_{n} \tag{1.3}
\end{equation*}
$$

with $\beta=1,2,4$ for the orthogonal, unitary, and simplectic Dyson ensembles, respectively.

For circular ensembles, we suggest here a qualitative interpretation of their physical meaning in terms of the time evolution, which is more relevant to the problems that will be considered later on. The eigenvalues of the operator $t \widehat{H}$ are simply the eigenenergies multiplied by the elapsed time, and they may be interpreted as phases accumulated by the states during the period of
time $t$. Exponentiation of $-i t \widehat{H}$ implies that the $n$-th eigenstate is multiplied by a phase factor $\exp \left(-i t E_{n}\right)$, or one can say that the real spectrum of $t \widehat{H}$ "winds" on the unit circle. Two energy eigenstates that are separated by an energy interval $2 \pi / t$ turn out to be at the same point on the unit circle. If the time interval $t$ is not too long, these states are far from each other and therefore do not experience a significant repulsion, typical of the neighboring energy eigenstates that result in the "rigidity" of the spectrum. However, for a time-dependent Hamiltonian there is a Fourier component of the interaction at the frequency $2 \pi / t$ which repels the states separated by the energy of the corresponding quantum as if they were neighboring. This provides us with an idea of how the eigenstates of the scattering matrix can form a"rigid" ensemble on the unit circle.

## Simple and Complex Multilevel Systems

A small number of multilevel problems have exact analytical solutions. The corresponding quantum systems usually have a certain, rather high, symmetry, either explicit or hidden in specific relations amongst the parameters of the unperturbed Hamiltonians $\widehat{H}_{0}$ and the perturbations $\widehat{V}$. This symmetry is actually the reason why the system dynamics, depending formally on a large number of the matrix elements of these operators, can still be described either in terms of a finite number of algebraic and special functions, or in terms of a few-fold integrals containing a limited number of combinations of these parameters. The symmetry underlying the exact solutions manifests itself in commutation relations between $\widehat{H}_{0}$, the perturbation $\widehat{V}$, and their commutators of higher orders, that altogether form a relatively simple operator algebra with the total number of elements much smaller than the total number of levels in the quantum system. In this case, the multilevel quantum system is in fact a multidimensional representation of a relatively small Lie group corresponding to the algebra.

One of the most illustrative examples is the magnetic multiplets of atomic levels with high angular momentum $J$ subjected to the action of magnetic fields - the population dynamics of the $2 J+1$ magnetic components is governed by the momentum operators $J_{+}, J_{-}$, and $J_{z}$ that form a closed threeelement algebra su(2). Another example is that of a harmonic oscillator in an external field, where the Hamiltonian $\widehat{a} \dagger \widehat{a}+1 / 2$, and the creation $\widehat{a} \dagger$ and annihilation operators $\widehat{a}$ also form a closed algebra.

However, this is not the case for the generic $N \times N$ matrices $\widehat{H}_{0}$ and $\widehat{V}$. Their higher-order commutators are usually linearly independent of the commutators of lower orders and the matrices themselves, and each new commutator yields a new linearly independent matrix until the total number of different matrices reaches the maximum possible number $N^{2}$. A statistical description seems the natural approach which allows one to make meaningful statements about the behavior of such systems. Of course, an important role is played by the quantitative relations between the commutators of different
orders. When a kind of hierarchic dependence exists, such that the sizes of the commutators rapidly decrease with increasing order of commutation, the same system may behave as a soluble system on a short time-scale but develop complex dynamics for a longer time domain. An example of such a system is a one dimensional chain of quantum states with disordered energies and the tight binding coupling, which demonstrates diffusive broadening of the population distribution over the states for short times but manifests the effect of quantum localization in the long time limit.

## Complexity and Chaos

One encounters complex quantum systems in physical problems of two different kinds. A problem of the first kind concerns a real physical ensemble of multilevel systems, that is with a set of quantum objects with the Hamiltonian of each object different from the Hamiltonian of another. One can find the ensemble average Hamiltonian and consider the differences between this average and Hamiltonians of particular systems as a random perturbation. We note that even a moderate perturbation of a multilevel system usually results in a large change of eigenstates and in an appreciable transformation of the energy eigenvalues distribution. Therefore many characteristics of a typical representative in the ensemble can be distinct from the mean characteristics over the ensemble, thus yielding a large dispersion.

A problem of the second kind relates to a single system with complex dynamics. In the classical limit such a system manifests chaotic behavior. Dealing with a classical chaotic system one never tries to follow the trajectory in phase space corresponding to a particular initial condition, but rather one takes a set of initial conditions and uses statistical characteristics such as Lyapunov exponents or Kolmogorov entropy for a description of the system dynamics. In this case the average over the initial conditions plays the role of the ensemble average.

We note however that in quantum mechanics one can usually find a natural representation for the Hamiltonian of a system of the second kind, where it has a simple structure and an explicit physical meaning. It might be for instance a product of physical quantities such as the coordinate, the electric field strength, the angular or kinetic momentum, etc. But in spite of this, a relatively small change in the Hamiltonian also results in a drastic change in the eigenstates and eigenvalues of the system. In this respect, a quantum system chaotic in the classical limit behaves like a system perturbed by a random matrix. The statement that all classically chaotic systems behave on the quantum level of consideration like the ensembles of random matrices is known as the Bohigas-Giannoni-Schmit conjecture.

This situation is typical for quantum systems of two or more spatial degrees of freedom. Indeed, if the unperturbed motion is separable, that is it can be represented as a combination of independent motions along each degree of freedom, the perturbation usually causes mixing of these motions and
destroys the separability of the problem. We can say that the perturbation destroys the quantum numbers, or in terms of classical mechanics, it destroys the integrals of motion. Another way to phrase the idea is to say that there are holonom constraints in the phase space of the separable classical system or in the Hilbert space of a quantum system. In classical mechanics these constraints are known as Kolmogorov-Arnold-Moser tori. The perturbation destroys the constraints and makes the system non-holonom.

In the following chapter we discuss a number of examples of complex quantum systems originating from different physical problems. Here we illustrate the analogy between a randomly perturbed quantum system and a system chaotic in the classical limit for one example, known as the Sinai billiard.


Fig. 1.1. The Sinai billiard and a random system with similar behavior.

We consider a particle in a two-dimensional potential well of infinite depth and with a square cross-section of side $L$. This motion of a particle in such a well is separable and the wavefunctions $\psi_{n, m}$ of the energy eigenstates are the products of the harmonic functions corresponding to each direction $\psi_{n, m}=(2 / L) \sin \left(k_{n} x\right) \sin \left(k_{m} y\right)$, where $k_{n}=2 \pi n / L$. We now place a scattering potential $V(x)$ at the corner of the well as shown in Fig. 1.1. In classical mechanics, for a round scatterer $V(x, y)=0$ at $x^{2}+y^{2}>a$ and $V(x, y)=V_{0} \rightarrow \infty$ at $x^{2}+y^{2}<a$ this system is known as a Sinai billiard where the motion of the particle becomes chaotic. Apparently it is also the case for a scatterer of an irregular (random) shape and for a large but finite $V_{0}$.

Let us now consider the problem quantum mechanically. The particle Hamiltonian reads

$$
\begin{equation*}
\hat{H}=\hat{H}_{0}+V(x, y)=\frac{1}{2 m} \hat{p}^{2}+V(x, y) \tag{1.4}
\end{equation*}
$$

and we assume that the potential $V(x, y)$ differs from zero in the region $S$ around the origin $x=y=0$. We take the eigenstates $\psi_{n, m}$ of the unperturbed Hamiltonian for the basis set of the problem and find the matrix elements of the interaction

$$
\begin{equation*}
V_{n, m}^{n^{\prime}, m^{\prime}}=\frac{4}{L} \int_{S} \sin \left(k_{n} x\right) \sin \left(k_{m} y\right) V(x, y) \sin \left(k_{n^{\prime}} x\right) \sin \left(k_{m^{\prime}} y\right) d x d y \tag{1.5}
\end{equation*}
$$

that turn out to be the Fourier transform of the potential. If $V(x, y)$ is a random function, the matrix elements are apparently random. But for a regular round scatterer they are definitely regular since they have the explicit form

$$
\begin{align*}
V_{n, m}^{n^{\prime}, m^{\prime}}= & \frac{4}{L} \int_{0}^{a} r d r \int_{0}^{\pi} d \phi \sin \left(k_{n} r \cos \phi\right) \sin \left(k_{m} r \sin \phi\right) \\
& V(x, y) \sin \left(k_{n^{\prime}} r \cos \phi\right) \sin \left(k_{m^{\prime}} \sin \phi\right) \\
= & \frac{2 V_{0}}{L^{2}} \sum_{s} c_{s} \int_{0}^{a} r d r J_{0}\left(\Delta k_{s} r\right)=\frac{2 V_{0}}{L^{2}} \sum_{s} c_{s} \frac{J_{1}\left(\Delta k_{s} a\right)}{\Delta k_{s} a}, \tag{1.6}
\end{align*}
$$

where $\Delta k_{s}= \pm \sqrt{\left(k_{n} \pm k_{n^{\prime}}\right)^{2}+\left(k_{m} \pm k_{m^{\prime}}\right)^{2}}, J_{0,1}(x)$ are the Bessel functions, and $c_{s}= \pm 1$.

We know, however, that both problems are equivalent in the classical limit. Therefore it is reasonable to assume that much of this similarity also persists in the quantum case. This implies that many of the predictions made for the systems perturbed by a random matrix will also be valid for a system perturbed by a non-random matrix if such perturbation destroys the separability of the problem and makes the motion chaotic, in agreement with the Bohigas-Giannoni-Schmit conjecture. In algebraic terms, the consecutive commutation of the regular perturbation matrix with the unperturbed Hamiltonian yields in this case the entire set of linearly independent matrices, exactly in the same way as commutation with a random matrix.

## Time Dependence and Spectral Properties

Our main concern will be to describe the population dynamics of complex quantum systems under the action of random perturbations. Many of the aspects of the time evolution have specific manifestations in the spectra of the ensemble averaged density of states $g(E)$, or more precisely, in the spectra of the total density of states projected either onto the initial state or onto a given state of the unperturbed system. This is usually the case when the dynamics are governed mainly by a specific unperturbed Hamiltonian, possessing some symmetries or corresponding to a separable motion of different degrees of freedom, while a random perturbation destroys this separability on a longer time-scale. The sharp edges of the state density $g(E)$ that are
already present for $E= \pm N\left\langle V^{2}\right\rangle$ in the Wigner distribution (1.2), a cusp on the dependence $g(E)$ at the point where two Wigner distributions merge, non-analytical dependencies $g(E)$ typical both for fractal structures and for systems corresponding to random walks in the classical limit, and some other spectral signatures, are usually associated with non-exponential decay processes. These processes are also known as coherent damping which usually corresponds to an intermediate asymptotic regime, that is to a relatively short time domain where the random perturbation has already started to affect the separability of the unperturbed system, but has not yet destroyed it completely.

Some other processes developing on a longer time-scale do not have signatures in the ensemble averaged density of states but manifest themselves in the statistics of the relative positions of the levels. These processes, known as recurrences and complete or fractional revivals, usually imply complete or partial returns of the population from different groups of quantum levels to the initial state. Interference of the returns from different groups of levels may result in an incomplete decay of the initially populated states which is one of the manifestations of so-called dynamical localization.

An interesting class of problems is represented by the case where the interaction which destroys the separability of the unperturbed motion changes significantly, by orders of magnitude, for different degrees of freedom. Therefore, in the course of time more and more degrees of freedom are getting involved in the dynamical process. The situation resembles the case known in classical mechanics as Arnold diffusion. For the corresponding quantum case, the recurrences and revivals play a very important role, systematically returning the population back to the initial state while still being unable to prevent this state from complete decay. The spectral density of the states accessible from the initial state is essentially a discontinuous and irregular function of the energy in this case. In the long time limit it corresponds to specific time dependences of the population of the initial state.

## Main Parameters and Quantitative Criterion of Complexity

The commutation relations between $\widehat{H}_{0}$ and $\widehat{V}$, generating at each order of commutation new linearly independent operators, cannot yet guarantee that the considered quantum object conforms to the model of a complex system. The quantitative parameters also play a crucial role, by indicating whether the action of the new operators can be considered as a small perturbation or whether the all-order perturbation theory is required for an adequate description. The main parameters that allow one to decide whether the system is complex are the average density of states $g$ of the unperturbed Hamiltonian and the mean squared strength of the perturbation $\left\langle V^{2}\right\rangle$.

The average density of states in the units where $\hbar=1$ has the dimensionality of time and represents the so-called Heisenberg time of return, a typical time when the return of the system to an initial state given by a superposition
of the eigenstates of $\widehat{H}_{0}$ has to be taken into account. The mean squared interaction enters the transition rate $2 \pi g\left\langle V^{2}\right\rangle$, suggested by the Fermi golden rule for the probability of transition per unit time from a given level to a group of levels with spectral density $g$, that are coupled to this level by an interaction characterized by the mean squared matrix element $\left\langle V^{2}\right\rangle$. Therefore the natural criterion of complexity reads $2 \pi g\left\langle V^{2}\right\rangle g \gg 1$ which implies that the probability of transition during the return time should be large.

Unless the special algebraic reasons emerging from the hidden symmetry of the problem result in the rapid decrease or vanishing of the commutators with the increase of the commutation order, the criterion $2 \pi g^{2}\left\langle V^{2}\right\rangle \gg 1$ also guarantees that the commutators can induce transitions. Since each new commutator is a linearly independent operator, the variety of linear combinations becomes more and more rich with the course of time. The interference of all of these transitions makes the dynamics of the system complex.

## Time-Dependent Perturbations and Control

A time-dependent perturbation enriches the system dynamics since each harmonic of the perturbation couples the levels separated by the energy of the corresponding quantum. For the periodic perturbation, the structure of the evolution operator is given by the Floquet theorem, which roughly speaking states that the amplitude $\bar{\psi}_{m}$ of each harmonic $m$ of the state vector can be considered as an independent variable in a Schrödinger equation with timeindependent parameters. The order of this matrix equation is the product of the dimension of the system Hilbert space and the total number of important harmonics. The order therefore approaches infinity when all harmonics are equally important. The only technical question arising for this approach is how to circumvent the ambiguity in relating the initial state of the system to the initial state of the harmonics.

For non-periodic perturbations, the time dependence can also be taken into account by introducing an additional degree of freedom which is now associated, not with a discrete index $m$ of $\bar{\psi}_{m}$, but with a continuous variable $\tau$ specifying the corresponding component $\bar{\psi}(\tau)$. In this new representation the Hamiltonians also experience a transformation by acquiring an additional operator term related to the new degree of freedom. Therefore the commutation relations between the unperturbed Hamiltonian and the perturbation may (and usually do) change, most strongly affecting "simple" multilevel systems which possess an underlying algebraic symmetry.

The time dependence may destroy the holonom constraints in a system including the energy conservation constraint, thus enriching the evolutionary capability of the latter. The question arises: can one employ this destruction and make the system evolve in an arbitrary predescribed way? In other words, can one force the system with a Hamiltonian $\widehat{H}_{0}$ to behave "a la carte", that is, as if it has another predetermined desired Hamiltonian $\widehat{H}_{d}$ ? The answer
to this question is positive for a system that conforms to the JurdjevicSussman theorem and corresponds to the compact algebra generated by the commutators of $\widehat{H}_{0}$ and $\widehat{V}(t)$. The main problem is how to construct an effective algorithm that would specify the required time dependence.

## Structure of the Text

The manner of presentation of the material serves our main objective, which is to invite the reader to consider different dynamical problems that belong to various domains of quantum physics from the point of view of complex quantum systems and attribute these problems to one of a number of different cases where exact and physically consistent conclusions can be drawn. In other words, the aim is to develop, as much as possible, the methods of solution rather than to calculate particular physical quantities for the particular physical problems. The latter will be simply impossible considering the vast number of domains where complex quantum systems can be encountered. Moreover, in some cases and for some physical properties, an alternative approach to the problem might be equally efficient, which does not however mean that the approach of the problem as a complex system is useless. Sometimes there may be additional insight into the reasons why the system behaves in such a way.

In Chap. 2, a number of examples are given in order to demonstrate the different physical domains where complex quantum systems can be encountered. Both the situation involving a physically existing ensemble of systems, and that involving a single system unknown in detail, are considered from the same standpoint, in other words, we do not specify whether a complex system means a single system with complex dynamics or an ensemble of systems, unless confusion is possible. Some of the examples are well known and widely investigated while others are rather new. Most of them are drawn from the fields of laser physics, physical chemistry, the physics of clusters, and quantum optics, where the dynamical aspects of multilevel quantum systems are the primary concern, whereas examples from solid state physics and nuclear physics - two domains where static aspects like quantum level statistics are explored by the method of random matrices and widely described in the literature - are presented hardly at all. In some examples we also show analogies with classical physics where the different modes of oscillations of a multidimensional system play the role of quantum states. The main purpose of this chapter is to show how one can describe the main features of the dynamical processes in terms of multilevel quantum systems.

In the next chapters we concentrate on the approaches to exactly soluble multilevel problems. Each chapter starts with a simple problem, usually wellknown, which serves the principal role of introducing the technical tools and for elucidation of the physical meaning and main parameters governing the dynamics. Technical complexity usually increases in the subsequent parts of the chapters. In these parts we have tried to give as many details as the
reader would need for following the narration without doing the calculations, or doing just a few of them.

In Chap. 3 we consider the two-level system and the level-band system, which are problems that can be solved exactly, without an ensemble average. We find the physical parameters responsible for the time evolution of such systems, and demonstrate the important role of coherence, which persists in the large multilevel systems and results in many interesting phenomena, such as quantum revivals, for instance. We discuss the role of the spectrum density and the spectrum edges, allow for the correlations in the matrix elements of interaction in the case of two levels interacting with a band resulting in specific spectral structures like the Fano profile, and analyze the situation when the interaction matrix elements may differ significantly, by orders of magnitude, resulting in non-exponential and incomplete decay of the level.

In Chap. 4 we consider the population dynamics in multilevel quantum systems perturbed by a time-independent random matrix. We illustrate the method of the ensemble average and demonstrate the technique, which can later be used for the description of more complicated problems. It is worth mentioning that in many respects, this technique is an analog of the diagram technique for the many body problem that has been successfully employed in solid state physics for the description of metals at low temperatures. We only employ the representation performed in the spirit of Grotrian diagrams which is more natural for atomic and molecular physics, where the unperturbed motion is depicted by levels (or dots) and the transitions are shown by lines. The main features of the diagram technique remain the same, including the idea of renormalization of the state energies, better known in atomic physics as "dressed states". We consider the role of the state degeneracy, the coherent damping associated with it, as well as the specific shapes of the spectral lines resulting from the coherence. We pay attention to the effects related to correlations in the perturbation matrices, when the rank of the perturbation matrix is much smaller compared to its order.

In Chap. 5 we consider the effect of a time-dependent perturbation on the population dynamics in "simple", exactly soluble quantum systems such as the two-level system and the level-band system. For a number of cases, the additional dimension associated with the time dependence can still result in an exactly soluble problem, although it complicates the algebraic relations amongst the Hamiltonian, the perturbation, and their commutators. The most well-known examples are the Landau-Zener problem and the DemkovOsherov problem. We also discuss the applicability of the semiclassical description also known as non-adiabatic transitions in time-dependent systems conforming to the Dykhne model, and discuss non-adiabatic transitions at the continuum edge.

In Chap. 6 we consider complex quantum systems subjected to a timedependent perturbation that also become exactly soluble following the ensemble average. We focus on the case where the system is perturbed by a
random matrix with a perturbation strength with an arbitrary time dependence. We demonstrate that the shape of the population distribution as a function of energy results from the interplay between the size of the perturbation and the spectral width of the time dependence. We also consider the harmonic and bi-harmonic perturbation of complex systems, as well as the bi-harmonic perturbation of a two-band system. The main purpose of this consideration is to illustrate how the additional degree of freedom, associated with the time dependence, manifests itself in the population distribution over the energy scale. We discuss the linear optical response, the delayed coherent polarization or optical echo, and the susceptibility of complex systems where one can see a difference between the systems corresponding to the unitary and the orthogonal ensembles. Controllability of complex quantum systems and several algorithms that offer the possibility to achieve complete control over the dynamics are discussed in this chapter as well.

In Chap. 7 we consider the situation where the levels are arranged in a one-dimensional relay-like chain and the interaction occurs only among the closest neighboring levels. We start with the problems that have exact analytical solutions, conforming to one of the low-dimensional algebras of operators. This concerns levels with linearly increasing energy and a constant coupling, the harmonic oscillator in the presence of dipole and Raman harmonic perturbations, angular momentum in an external magnetic field, and a few other examples where the exact solution in terms of hypergeometric functions exists. For the case of smooth variation of the parameters, we describe an approach in the spirit of the WKB approximation. We also consider the one-dimensional chain of irregularly detuned levels, the so-called Lloyd chain, whose spectrum and localized wavefunctions have been widely explored, but whose population dynamics have been much less addressed. In this context we introduce the field-theory methods (also convenient for some other applications) that allow one to calculate ensemble averaged determinants of random matrices, representing them as functional integrals over regular and Grassmann variables.

In Chap. 8 we consider the dynamics of composed complex quantum systems where the perturbation matrix has a block structure. We discuss the diffusion regime in the relay-like structures for various dependencies of the transition rate on the block number, that allow an exact description in terms of the hypergeometric functions, and show the important role of the correlations in the perturbation matrix elements. We also concentrate on the dynamics of multilevel systems at long times. We illustrate the dominating role of the quantum interference effects on the population dynamics. In particular, we calculate the probability for a complex system to stay in the initial state at asymptotically long time. This characteristic is complementary to the property of quantum localization in disordered systems. We discuss the manifestation of quantum complexity in the density of states and in the dipole absorption profiles along with a relatively simple way to calculate these pro-
files for the case when the spectral structures typical of the unperturbed Hamiltonian are still not completely washed out by the random perturbation. We demonstrate the spectral manifestation of the hierarchical structure of the system, originating from the fractal nature of an object and the transformations that this spectrum experiences under a random perturbation.

There is no literature references throughout the text, apart from the names of the authors traditionally associated with the well-known results of methodical importance. One of the reasons for this is the fact that many similar or sometimes even identical results have been independently obtained in different domains of physics and are known under different names. For instance, the dipole susceptibility of highly exited molecules has almost the same structure as conductance in metals, while the tunneling transparency of disordered solids has much in common with the cross-section of multiphoton ionization of atoms. The conductivity of narrow-band solids has exactly the same mathematical description as the multiphoton off-resonant excitation of ro-vibronic states of diatomic molecules. The interaction of resonant atoms in a photonic band gap is similar to the Raman transitions between two vibrational states of a molecule. The Anderson localization turns out to be analogous to the Fermi acceleration of the quantized kicked rotor, and the Breit-Wigner distribution in nuclear physics is identical to the Lorentzian profiles of optical line shapes, and also similar to Feshbach resonances for molecular collisions. The idea of renormalization, widely employed in particle physics, practically coincides with the concept of dressed states for the atoms in strong laser fields, etc. In such a situation, the references to different physical domains throughout the text would be confusing. That is why the review of the literature is performed in a separate, and last, Chap. 9, although even there the list of literature cannot be complete and in fact is far from being exhaustive. A relatively short review is presented in the form of problems which the reader can consider with the help of the approaches developed and compare the formulations of the problems and the results with the original papers. More detailed information is available in the monographs listed in accordance with the distribution of topics among the chapters.

However, the reader is assumed to know the main concepts and methods of quantum mechanics and mathematical methods of physics. For the fundamentals of quantum mechanics one can consult the classical series of textbooks by L. D. Landau and E. M. Lifshits as well as Quantum mechanics by C. Cohen-Tanoudji, B. Diu, and F. Laloe [49]. The main mathematical tools can be found in the textbook by J. Mathews and R. L. Walker [134] Mathematical methods of physics. For more special aspects of the mathematical technique including the algebra and group theory, the books Mathematical methods of classical mechanics by V. I. Arnold, A. Weinstein, and K. Vogtmann [9], Theory of group representations and applications by A. Barut and R. Raczka [12], Integrable systems of classical mechanics and Lie algebras by A. M. Perelomov [150], and the article Universal integrals of motion
and universal invariants of quantum systems by V. V. Dodonov [61] can be recommended.

## 2 Examples of Complex Systems

Molecular and laser physics offer many examples of complex quantum systems that demonstrate the importance of characteristics in disordered spectra other than the density of states or density of states correlations. Here we briefly discuss some of these systems that have been studied intensively because of their fundamental interest or practical importance.

We start the consideration with problems related to the domain usually known as the interaction of radiation with matter which is dedicated to atoms and molecules in intense resonant laser fields. In these types of problems, one can usually separate one or a few optically active degrees of freedom that interact with the laser field directly, while the rest of the degrees of freedom are not dipole active, but simply coupled to the optically active ones by intra-molecular or intra-atomic interactions. The presence of the latter make the system complex. Highly exited atomic states that belong to Rydberg series have a rich spectrum and also form a complex system, being perturbed by a strong laser or RF field. However the multilevel system should not necessarily correspond to highly exited states. Polyatomic molecules, even in the ground vibrational state, posses a rather rich rotational spectrum, which can be strongly affected by an external field. Additional complexity is associated with the presence of multiphoton resonances that correspond to the simultaneous absorption of several laser quanta.

Chemical reactions offer a large number of examples of complex quantum systems. Each molecule is already a complex object in itself, whereas two colliding molecules, even in the absence of chemical reaction, form a complex compound with a rich spectrum of vibrational states, the population of which changes significantly during the collision. An additional source of complexity is the non-adiabatic interaction of the electronic and vibrational degrees of freedom inside a single vibrationally exited molecule which results in intramolecular energy conversion.

The situation becomes even more involved in the case when the molecules are highly exited electronically, approaching molecular Rydberg states. In such a situation, the adiabatic separation of the fast electronic and slow nuclear motions fails, since the period of rotation of an electron along a Rydberg orbit becomes comparable to, and sometimes even longer than, the vibrational period. The electronic and vibrational states get strongly entangled
and form a common complex spectrum. The collision of a Rydberg molecule with another molecule may result in collisional ionization which adds a continuous component to the system spectrum. An even more complicated situation arises when we consider an ensemble of many interacting Rydberg atoms. Excitation exchange among the atoms results in the motion of populations over the volume occupied by the gas. The spectrum now becomes very rich and involved, corresponding to the mobility of different exitons over the volume.

Another possibility to form complex systems is to put many identical atoms together, such that they form an atomic cluster. Even in the case when each atom has only two quantum states, the spectrum of a compound containing $n$ atoms already has $N=2^{n}$ quantum levels and becomes a complex system in the presence of even a very weak interatomic interaction. This spectrum becomes completely disordered if we take into account the interaction of the electrons with atomic vibrations or phonons. This quantum complexity manifests itself in the thermodynamical as well as in the optical properties of the clusters. The complex spectra of clusters are not necessarily formed as the result of atomic forces. The interaction can be mediated by the electromagnetic field. In this case one can see the cooperative interaction of the atoms with the external electromagnetic field.

Quantum complexity is not only typical of atoms and molecules. It has parallels with multidimensional classical systems such as semiclassical translational motion of ions in ionic traps, nonlinear interaction of photons in a resonator, and mobility of electrons on surfaces. In this chapter we show how the equation of motion can be written in a similar form for all of the aforesaid examples. All of the systems, including classical systems, conform to one of many possible models of complex quantum systems.

### 2.1 Molecules and Atoms in Laser Fields

### 2.1.1 Laser Breaking of a Weakly Bonded Complex

The difference in the vibrational frequencies of a given molecule containing different isotopes is exploited by the technique of laser isotope separation. Let us consider in this context a coordinated compound that is a simple molecule, $\mathrm{IsO}_{2}$ for example, which is surrounded by a number of relatively weakly bound molecular complexes or ligands as shown in Fig. 2.1. Here Is stands for an element that has variable isotopic modifications. The central molecule is subjected to the action of the infrared laser which is in resonance with a single dipole-active vibrational mode. The ligands may both alter the frequency of the vibration, and quench the vibrational excitation of the central molecule. The energy of one vibrational quantum of the excited mode is usually sufficient to destroy the weakly bound ligands, and therefore the transfer of the vibrational energy from the molecule to the surrounding complexes immediately results in the dissociation of the compound.


Fig. 2.1. A weakly bound dioxide compound. The dipole-active vibrations of the oxygen atoms are much faster than the chaotic multidimensional motion of the molecular groups surrounding the dioxide.

How can one formulate this problem mathematically? We use two indices to describe the quantum state of the system - the index $n$ which enumerates the vibrational state $\psi_{n}(q)$ of the mode $q$ of the central molecule excited by the radiation, and the index $m$ which enumerates the energy eigenstate of the entire system of the molecular compound provided the vibrational state $n$ of the excited mode is fixed. We denote by $\Phi_{n, m}(Q)$ the corresponding wavefunctions of the compound, where $Q$ stands for all of the set of $N$ coordinates of the compound other than $q$. We note that the Hamiltonian of motion for the surrounding molecules depends on the vibrational state of the central molecule, since the vibrationally excited central molecule has both a different size and different force constants for the weak bonds compared with the ground vibrational state of the central molecule. We therefore write the Hamiltonian of the entire compound in the form

$$
\begin{equation*}
\hat{H}=n \hbar \omega_{I s O_{2}} \delta_{n, n^{\prime}} \delta_{m, m^{\prime}}+H_{m, m^{\prime}}(n) \delta_{n, n^{\prime}} \tag{2.1}
\end{equation*}
$$

where we still neglect the influence of the energy transfer from the central molecule to the periphery, and the excitation by the laser field. The eigenenergies

$$
\begin{equation*}
E_{n, m}=\hbar \omega n+\Delta_{n, m} \tag{2.2}
\end{equation*}
$$

correspond to the eigenfunctions $\left|\Psi_{n, m}\right\rangle=\left|\psi_{n}\right\rangle\left|\Phi_{n, m}\right\rangle$.
Now let us allow for the external resonant laser field which induces transitions among the states of the system. The probability amplitudes of the transitions are given by the product of the electric field strength $E=\mathcal{E} \cos \omega t$ and the matrix elements of the dipole moment operator $d_{n, m}^{n^{\prime}, m^{\prime}}=\langle n, m| e \hat{q}\left|n^{\prime}, m^{\prime}\right\rangle$.

The latter is the product of two parts. The first part originates from the excited vibrational mode of the central molecule and is given by the dipole matrix element of the fundamental transition, whereas the second part arises from the overlap integrals of the pairs of wavefunctions of the surroundings, corresponding to the central molecule in the ground vibrational state and in the first excited vibrational state, respectively. It reads

$$
\begin{align*}
d_{n, m}^{n^{\prime}, m^{\prime}} & =\int \psi_{n}(q) e q \psi_{n \pm 1}(q) d q \int \Phi_{n, m}(Q) \Phi_{n \pm 1, m^{\prime}}(Q) d^{N} Q \\
& =d \sqrt{n} \int \Phi_{n, m}(Q) \Phi_{n \pm 1, m^{\prime}}(Q) d^{N} Q \tag{2.3}
\end{align*}
$$

Here we have employed the harmonic oscillator model for the vibrational modes of the central molecule which gives the selection rules $n-n^{\prime}= \pm 1$ and have denoted by $e$ the effective charge of the mode $q$.

Now we assume that the motion of the surroundings is chaotic and therefore a small change in the Hamiltonian of this motion associated with the excitation of the central molecule results in a big change of the eigenfunctions $\Phi_{n, m}(Q)$. In other words the vibrational excitation of the mode $q$ results in a rotation of the ligand energy eigenstates basis in Hilbert space by large angles around many axes. Therefore for an ensemble of molecules, the statistical approach is the natural way to describe the quantum dynamics, and the normal distribution is the logical assumption for the distribution functions of the overlap integrals $\left\langle\Phi_{n, m} \mid \Phi_{n^{\prime}, m^{\prime}}\right\rangle$ in (2.3).

One can write the Schrödinger equation for the amplitudes $\Psi_{n, m}$ of the states $\left|\Psi_{n, m}\right\rangle$ in the form

$$
\begin{equation*}
i \hbar \dot{\Psi}_{n, m}=\left(\hbar \omega_{I s O_{2}}+\Delta_{n, m}\right) \Psi_{n, m}+\cos \omega t \sum_{m^{\prime}} V_{n, m}^{n \pm 1, m^{\prime}} \Psi_{n \pm 1, m^{\prime}} \tag{2.4}
\end{equation*}
$$

where $V_{n, m}^{n \pm 1, m^{\prime}}=\mathcal{E} d_{n, m}^{n \pm 1, m^{\prime}}$.
If we now take the resonant approximation, that is we take the amplitudes in the form

$$
\begin{equation*}
\Psi_{n, m}=\psi_{n, m} \exp \{i n \omega t\} \tag{2.5}
\end{equation*}
$$

substitute (2.5) into (2.4), and neglect the rapidly oscillating terms, we arrive at

$$
\begin{align*}
i \hbar \dot{\psi}_{n, m}= & \left(\omega_{I s O_{2}}-\omega+\Delta_{n, m}\right) \psi_{n, m} \\
& +\sum_{m^{\prime}}\left(V_{n, m}^{n+1, m^{\prime}} \psi_{n+1, m^{\prime}}+V_{n, m}^{n-1, m^{\prime}} \psi_{n-1, m^{\prime}}\right) \tag{2.6}
\end{align*}
$$

This is the Schrödinger equation for the slow amplitudes $\psi_{n, m}$. The Hamiltonian of this equation contains two components which we shall consider as random: $\Delta_{n, m}$ gives us the random deviations of the level energies from their harmonic positions $\hbar n \omega$, while $V_{n, m}^{n \pm 1, m^{\prime}}$ allows for the random coupling of
these states. The problem now is to find some average characteristics of these matrices which governs the population dynamics of the system. That would allow us to identify the common properties of all similar systems, provided such common properties exist.

### 2.1.2 Laser-Induced Electronic Transitions in Molecules

There is a number of physical systems that are similar to the one just considered. Thus far we were considering the excitation of a vibrational mode of the central molecule and the effect of this vibration on the weak bonds in the compound. The crucial step there was the adiabatic separation of the rapid vibrations and the slow motion of the surrounding molecules. This is completely analogous to the adiabatic separation of the electronic and vibrational motions in molecules, and therefore we face the same situation considering the electronic excitation of polyatomic molecules in laser fields.

This is illustrated in Fig. 2.2. The laser field induces transitions among


Fig. 2.2. Electronic transition in a molecule.
the manifolds of vibrational states corresponding to two different electronic terms. Transition amplitudes are proportional to the Frank-Condon factors, that is the overlap integrals $\left\langle\Phi_{m}^{0} \mid \Phi_{m^{\prime}}^{1}\right\rangle$ of the vibrational wavefunctions $\Phi_{m}^{0}$ in the ground electronic term, and the vibrational wavefunctions $\Phi_{m^{\prime}}^{1}$ in the excited term. The corresponding Schrödinger equation for slow amplitudes $\psi_{n, m}$ of the vibrational states $\left|\Phi_{n, m}\right\rangle$ of the $n$-th electronic term $(n=0,1)$ reads

$$
\begin{align*}
& i \hbar \dot{\psi}_{1, m}=\left(E_{1, m}-\hbar \omega\right) \psi_{1, m}+\sum_{m^{\prime}} V_{1, m}^{0, m^{\prime}} \psi_{0, m^{\prime}} \\
& i \hbar \dot{\psi}_{0, m}=E_{0, m} \psi_{0, m}+\sum_{m^{\prime}} V_{0, m}^{1, m^{\prime}} \psi_{1, m^{\prime}} \tag{2.7}
\end{align*}
$$

where $V_{n, m}^{n^{\prime}, m^{\prime}}$ now stands for the probability amplitudes of the laser-induced transitions.

The nuclear motion of vibrationally excited polyatomic molecules is usually chaotic. This implies that the vibrational wavefunction changes drastically and in an irregular way when we move from a given energy eigenstate to a neighboring one. This is even more the case for the vibrational functions corresponding to different electronic terms, which implies that the overlap integrals $\left\langle\Phi_{m}^{0} \mid \Phi_{m^{\prime}}^{1}\right\rangle$ change completely with any change of the indexes $m$ and $m^{\prime}$. One, therefore, may consider the matrix element of the dipole transition

$$
\begin{equation*}
V_{0, m}^{1, m^{\prime}}=\mathcal{E} d_{0, m}^{1, m^{\prime}}=\left\langle\phi_{0}\right| \mathcal{E} \hat{d}\left|\phi_{1}\right\rangle\left\langle\Phi_{m}^{0} \mid \Phi_{m^{\prime}}^{1}\right\rangle \tag{2.8}
\end{equation*}
$$

as a random value. Here $\phi$ denotes the electronic wavefunctions, and $\Phi$ the vibrational wavefunctions dependent on all the nuclear coordinates.

### 2.1.3 Vibrational Excitation of Polyatomic Molecules

We encounter the same situation for the vibrational motion of a polyatomic molecule excited by an infrared laser field resonant to one of its vibrational modes. In this case the Schrödinger equation formally coincides with (2.6) where $n$ denotes the number of absorbed laser quanta. Let us consider it in more detail. The anharmonicity of vibrations makes the motion of the polyatomic molecule chaotic, and hence the matrix elements of the dipole moment acquire a random character. We illustrate this situation in the most simple example of a three-fold degenerate mode of a symmetric polyatomic molecule $X Y_{4}$ or $X Y_{6}$. In the harmonic approximation excited vibrational levels of the mode are multiply degenerate. To each eigenstate of the degenerate mode $\left|I_{1}, I_{2}, I_{3}\right\rangle$ one can attribute three quantum numbers $I_{1}, I_{2}$, and $I_{3}$, that correspond to numbers of quanta in each of three vibrational degrees of freedom comprising the mode. The energy of the state reads $E_{I_{1}, I_{2}, I_{3}}=\hbar \omega_{v i b}\left(I_{1}+I_{2}+I_{3}\right)$, where the vibrational frequency $\omega_{v i b}$ coincides with the frequency of the external laser field $\omega$. For light absorption three possible dipole transitions from the state $\left|I_{1}, I_{2}, I_{3}\right\rangle$ are allowed by the selection rules. They correspond to an increase by one of the number of vibrational quanta in each degree of freedom:

$$
\left|I_{1}, I_{2}, I_{3}\right\rangle \rightarrow\left\{\begin{array}{l}
\rightarrow\left|I_{1}+1, I_{2}, I_{3}\right\rangle  \tag{2.9}\\
\rightarrow\left|I_{1}, I_{2}+1, I_{3}\right\rangle \\
\rightarrow\left|I_{1}, I_{2}, I_{3}+1\right\rangle
\end{array}\right.
$$

We now take into account the anharmonicity of the vibrations. In a symmetric, three-dimensional oscillator the anharmonicity consists of three parts.

The first part $\hat{H}_{1}$ is the regular, one-dimensional anharmonicity which shifts the positions of levels of each degree of freedom independently:

$$
\begin{equation*}
\hat{H}_{1}\left|I_{1}, I_{2}, I_{3}\right\rangle=\alpha\left(I_{1}^{2}+I_{2}^{2}+I_{3}^{2}\right)\left|I_{1}, I_{2}, I_{3}\right\rangle \tag{2.10}
\end{equation*}
$$

where by $\alpha$ we denote the corresponding anharmonicity constant. The second part is the cross-anharmonicity. It allows for the change in frequency of one degree of freedom following the excitation of another. We denote the corresponding constant $\beta$ and write the Hamiltonian in the form

$$
\begin{equation*}
\hat{H}_{2}\left|I_{1}, I_{2}, I_{3}\right\rangle=\beta\left(I_{1} I_{2}+I_{2} I_{3}+I_{3} I_{1}\right)\left|I_{1}, I_{2}, I_{3}\right\rangle \tag{2.11}
\end{equation*}
$$

Neither the first nor the second part allows for energy transfer between the degrees of freedom of the degenerate mode. They do not destroy the separability of the problem and each degree of freedom may thus far be considered as dynamically independent. The third part of the anharmonicity, however, is responsible for the mode mixing. The Hamiltonian $\hat{H}_{3}$ corresponding to this part is non-diagonal in the representation of the quantum numbers $I_{1}, I_{2}, I_{3}$, and mixes the states

$$
\begin{equation*}
\left|I_{1}, I_{2}, I_{3}\right\rangle \leftrightarrow\left|I_{1}-2, I_{2}+2, I_{3}\right\rangle \leftrightarrow\left|I_{1}-2, I_{2}, I_{3}+2\right\rangle \leftrightarrow \text { etc. } \tag{2.12}
\end{equation*}
$$

by transferring pairs of quanta from one degree of freedom to another. It does not commute with $\hat{H}_{1}$ and $\hat{H}_{2}$. One can write down this part of the anharmonicity Hamiltonian in terms of the creation $\hat{a}_{j}^{\dagger}$ and annihilation $\hat{a}_{j}$ operators of the vibration of the $j$-th degree of freedom:

$$
\begin{equation*}
\hat{H}_{3}=\gamma\left(\hat{a}_{1} \hat{a}_{1} \hat{a}_{2}^{\dagger} a_{2}^{\dagger}+\hat{a}_{1} \hat{a}_{1} \hat{a}_{3}^{\dagger} a_{3}^{\dagger}+\hat{a}_{2} \hat{a}_{2} \hat{a}_{3}^{\dagger} a_{3}^{\dagger}+\text { c.c. }\right) \tag{2.13}
\end{equation*}
$$

hereafter c.c. denotes the complex conjugate expression. Note, that the third part of the anharmonic Hamiltonian commutes with the total number $I=$ $I_{1}+I_{2}+I_{3}$ of the quanta in the mode, and leaves this quantity intact.

Diagonalization of the sum of the three Hamiltonians $\hat{H}_{1}+\hat{H}_{2}+\hat{H}_{3}$ at fixed $I$ gives the eigenstates $|m I\rangle$ of the three-fold degenerate mode. Here $m$ enumerates the levels corresponding to a given $I$. The many-fold degeneracy of the state is now raised, as shown in Fig. 2.3. Usually the results of diagonalization are very sensitive to the values of the anharmonicity constants $\alpha, \beta$, and $\gamma$, and both the eigenvalues $E_{I, m}$ and the eigenstates $|m I\rangle$ of the Hamiltonian change drastically with a relatively small change in each of the parameters. This is typical of chaotic systems and implies that if we cast the eigenstates $|I, m\rangle$ of the three-dimensional anharmonic oscillator in terms of linear combinations of harmonic states $\left|I_{1}\right\rangle,\left|I_{2}\right\rangle$, and $\left|I_{3}\right\rangle$ :

$$
\begin{equation*}
|I, m\rangle=\sum_{I_{2}, I_{3}} c_{I_{2}, I_{3}}^{I, m}\left|I-I_{2}-I_{3}, I_{2}, I_{3}\right\rangle \tag{2.14}
\end{equation*}
$$

the expansion coefficients $c_{I_{2}, I_{3}}^{I, m}$ will be complex functions of $I_{1}, I_{2}$, and $m$ that strongly depend on $\alpha, \beta$, and $\gamma$. Therefore the matrix of dipole transitions among the states $|I m\rangle$


Fig. 2.3. Anharmonic splitting of a three-fold degenerate mode.

$$
\begin{align*}
\langle I, m| \hat{d}\left|I+1, m^{\prime}\right\rangle= & \sum_{I_{2}, I_{3}}\left(c_{I_{2}, I_{3}}^{I, m}\left\langle I_{1}\right| \hat{d}\left|I_{1}+1\right\rangle c_{I_{2}, I_{3}}^{I+1, m^{\prime}}\right. \\
& +c_{I_{2}, I_{3}}^{I, m}\left\langle I_{2}\right| \hat{d}\left|I_{2}+1\right\rangle c_{I_{2}+1, I_{3}}^{I+1, m^{\prime}} \\
& \left.+c_{I_{2}, I_{3}}^{I, m}\left\langle I_{3}\right| \hat{d}\left|I_{3}+1\right\rangle c_{I_{2}, I_{3}+1}^{I+1, m^{\prime}}\right) \tag{2.15}
\end{align*}
$$

is an irregular matrix constrained only by the condition $I-I^{\prime}= \pm 1$. If we replace it by a random matrix we arrive at the Schrödinger equation

$$
\begin{align*}
i \hbar \dot{\psi}_{I, m}= & \left(E_{I, m}-I \hbar \omega\right) \psi_{I, m} \\
& +\sum_{m^{\prime}}\left(V_{I, m}^{I+1, m^{\prime}} \psi_{I+1, m^{\prime}}+V_{I, m}^{I-1, m^{\prime}} \psi_{I-1, m^{\prime}}\right), \tag{2.16}
\end{align*}
$$

which has the same structure as (2.6).

### 2.1.4 Transitions Among Levels with Fine Structure

We start with the fine structure of the rotational states of spherical-top symmetric polyatomic molecules which also conform with the random matrix model. Spherical-top molecules have three equal components of the inertia tensor, and therefore the rotational states are degenerate with respect to the quantum number $K$, which gives the projection of the vector of total angular momentum $\boldsymbol{L}$ to the molecular axis. This may be interpreted that the molecule has the same rotational energy regardless of the orientation of the rotation axis in the reference system associated with the atomic positions.

This, however, is not the case for a non-rigid top, if its mechanical elasticity depends on the orientation of the rotation axis. For a given angular velocity $\boldsymbol{\Omega}$ the rotational energy of the top is dependent on the orientation of the rotation axis, and the angular momentum $L$ is a nonlinear vector function of $\boldsymbol{\Omega}$. The direction of $\boldsymbol{L}(\boldsymbol{\Omega})$ may deviate from the direction of $\boldsymbol{\Omega}$ and
this disalignment varies with the relative orientation of the molecule. For a classical top the conservation of energy and angular momentum implies that the magnitude and the direction of $\boldsymbol{\Omega}$ changes in time, which also implies that the orientation of the molecule with respect to the angular momentum vector changes. The time evolution of the orientation of a non-rigid molecule has an irregular and chaotic character. At the quantum level of consideration this means that the quantum number $K$ does not correspond to a certain energy eigenstate, that is the $K$-degeneracy is raised, and quantum numbers $K$ are destroyed. We mark the quantum states $|n\rangle$ of the resulting multiplets with the index $n$, which now has nothing to do with the angular momentum projection. The number of these states is given by the number $2 L+1$ of different $K$-components and can be large for rapidly rotating molecules.

Consider now the dynamics of excitation of vibrational-rotational transitions in spherical-top molecules by an external laser field. The molecule in the ground vibrational state can be considered as rigid, and the rotational levels are $K$-degenerate. In an excited vibrational state it looses the rigidity and the $K$-degeneracy is raised. We denote by $\Delta_{n}$ the displacement of the states $|n\rangle$ from their unperturbed positions and write the Schrödinger equations for the slow amplitudes in the form

$$
\begin{align*}
i \hbar \dot{\psi}_{K} & =\sum_{n} V_{K, n} \psi_{n} \\
i \hbar \dot{\psi}_{n} & =\Delta_{n} \psi_{n}+\sum_{K} V_{n, K} \psi_{K} \tag{2.17}
\end{align*}
$$

where $\psi_{K}$ correspond to the ground vibrational state, and $\psi_{n}$ denotes the amplitudes of the rotational energy eigenstates of the vibrationally excited molecule. The total momentum $L$ and its projection $M$ are supposed to be conserved in the course of molecular rotations. Note that the displacements $\Delta_{n}$ may also include the detuning of the laser field from the frequency of the unperturbed transition.

For a one-photon process the transition probability amplitudes $V_{n, K}$ are given by the product of the laser field amplitude $\mathcal{E}$, the dipole moment of vibrational transition $d_{v i b}$, and the scalar product of states $\langle n \mid K\rangle$. Chaotic motion of the non-rigid top suggests the random matrix model for the latter. For a multiphoton excitation we have to replace the transition amplitudes $V_{n, K}=\mathcal{E} d_{v i b}\langle n \mid K\rangle$ by the composite matrix elements

$$
V_{n, K}=\langle n \mid K\rangle \mathcal{E} d_{1} \prod_{l=2}^{\substack{\text { number of } \\ \text { absorbed photons }}}\left\{\frac{\mathcal{E} d_{l}}{E_{l}-l \hbar \omega}\right\},
$$

where $E_{l}$ represents the energies of the intermediate states, and $d_{l}$ the dipole moments for the transitions $l-1 \rightarrow l$.

### 2.1.5 Excitation of Rydberg States in Atoms

Degeneracy of states is also present in atoms. However the optical transitions in atoms usually occur among states with relatively small angular momenta, where the number of the magnetic components is relatively small. The population dynamics in atoms can therefore be studied analytically or numerically, where an approach based on statistical ideas is not really needed. Nevertheless they are useful when we consider multiphoton transitions to highly excited Rydberg states. Indeed, multiphoton excitation can be associated with a large acquisition of angular momentum, and thus highly excited states with large angular momenta have rich multiplet structures resulting from the degeneracy of the magnetic quantum numbers $M$. In the presence of electric and magnetic fields the joint action of Stark and Zeeman effects raises the degeneracy and destroys the quantum numbers $M$. The situation is identical to that of spherical-top molecules; many components of the multiplets simultaneously fall into a multiphoton resonance with the external laser field, and no selection rules are left. One may use (2.17) with the composed matrix elements similar to those given in (2.18).

Another example of a complex multilevel system is that of Rydberg atoms in an ultra-high-frequency radio field. One can formally describe the system with the help of (2.16) where $I$ denotes the principal quantum number, and $m$ stands for the magnetic quantum number. However, the matrix elements $V_{n, l}^{n^{\prime}, l^{\prime}}$ cannot a priori be taken as random due to the presence of selection rules. Nevertheless, since the classical motion of an electron in the Coulomb and radio-frequency fields is chaotic, the corresponding quantum problem


Fig. 2.4. Rydberg atoms versus polyatomic molecules in laser fields. The spectral density of the vibrational states of a polyatomic molecule rapidly increases with energy (a). The laser field induces upward transitions. In the compound system of the molecule+quantized field (b) the transitions occur among the states of the same energy. An atom excited by a laser field (d) to a Rydberg state in the presence of an RF field performs a series of transitions (f) towards the ionization limit (e)
may manifest some properties of a multilevel system perturbed by a random matrix, although the separation of the Hamiltonian into a regular part and a random part may be different from the separation of the Hamiltonian of the Coulomb problem and the RF-field dependent perturbation. Such a point of view is complementary to the traditional approach to the problem which relies on the numerical investigation of quantum maps. In Fig. 2.4 we compare and contrast Rydberg atoms and molecules in an oscillating external field.

### 2.1.6 Competition of Multiphoton Processes of Different Orders

For molecules in strong laser fields one can face competition between multiphoton processes of different orders. Indeed, polyatomic molecules posses rich vibrational spectra. The detuning of the lowest excited level of the vibrational mode from one-photon resonance can be equal or close to the detunings of higher levels from the higher-order multiphoton resonances ( Fig. 2.5(a)). The larger the order of the resonance $n$ the higher the density of states $g_{n}$ in the vicinity of the resonant energy $n \hbar \omega$. The excitation dynamics of such a system result from the competition between the increasing density of states and decreasing size of the composed matrix elements (2.18) of the multiphoton transitions. We illustrate such a situation in Fig. 2.5.

In the resonant approximation this system is equivalent to the system of interacting states shown in Fig. 2.5(b), where a large matrix element $V_{1}$ connects the starting state $|0\rangle$ with a small number of levels, and a weaker interaction $V_{2}$ corresponds to a higher density of states $g_{2}$, etc. The most dense spectrum is hardly coupled to the starting state.


Fig. 2.5. (a) Competition between multiphoton processes of different orders. (b) Density of resonances as a function of excitation energy. (c) Level scheme for resonances of different orders. The increasing density of resonances $g_{n}$ compensates for the decreasing coupling $V_{n}$.

The Schrödinger equation for such a system reads

$$
\begin{align*}
i \hbar \dot{\psi}_{0} & =\sum_{m} V_{0, m} \psi_{m} \\
i \hbar \dot{\psi}_{m} & =\Delta_{m} \psi_{m}+V_{m, 0} \psi_{0} \tag{2.19}
\end{align*}
$$

where $\psi_{0}$ denotes the amplitude of the ground state of the system, and the excited states with the amplitudes $\psi_{m}$ are detuned from the resonant positions of one- or multiphoton resonances by the energy $\Delta_{m}$.

One of the possibilities to tackle the problem is to model the interaction and detunings by random values. This model should however allow for large differences in the typical matrix elements of interactions, that correspond to different orders of resonances. We note the principal difference between this problem and the problem of a level decaying to a continuum: two close states $|m\rangle$, and $\left|m^{\prime}\right\rangle$ in the system described by (2.17) which have a small difference in the detunings $\Delta_{m}-\Delta_{m^{\prime}}$ may possess transition amplitudes $V_{m}$ and $V_{m^{\prime}}$, different by orders of magnitude.

### 2.2 Collisions and Reactions of Molecules

Interaction of atoms and molecules with laser radiation is not the only domain where consideration of the population dynamics is needed. A great number of similar problems arise in chemistry, in the theory of atomic and molecular collisions, in the theory of Rydberg molecules, associative Penning ionization, etc. In this section we consider some of these examples.

### 2.2.1 Collisional Redistribution of Energy

Let us first consider the simplest example of a polyatomic molecule colliding with a simple, namely structureless, particle. We assume that the trajectories of the colliding particles are given and concentrate only on the change of the internal state, that is the vibrational state of the polyatomic molecule. Since vibrationally excited polyatomic molecules move chaotically, we also assume that the interaction Hamiltonian $\hat{H}_{\text {int }}$ in the representation of the vibrational energy eigenstates $|n\rangle$ is complex, and is given by a random matrix $\hat{H}_{\text {int }}=G(t) V_{n, n^{\prime}}$, where the size of the interaction $G(t)$ depends on the distance between the colliding particles, and therefore changes with time. We write the Schrödinger equation for the polyatomic molecule

$$
\begin{equation*}
i \hbar \dot{\psi}_{n}=E_{n} \psi_{n}+G(t) \sum_{m} V_{n}^{m} \psi_{m} \tag{2.20}
\end{equation*}
$$

which differs from (2.16), since in the interaction term it contains the time dependence explicitly. The redistribution of the population among the vibrational states of the molecule (see Fig. 2.6) depends on the typical size $V$ and


Fig. 2.6. Redistribution of the vibrational state populations in the course of the collision of a polyatomic molecule with a buffer gas molecule.
the time dependence $G(t)$ of the interaction. The solution of (2.16) allows us to find the proper combination of $V$ and $G(t)$ that governs the probability of the energy exchange.

Collisions among vibrationally excited polyatomic molecules also result in vibrational energy exchange. If the distance between two colliding molecules at the position of the closest approach exceeds their sizes, we can assume that the translational and the vibrational degrees of freedom are decoupled from each other and take straight lines for the trajectories of the centers of mass of the molecules. The main contribution to the energy transfer then results from the dipole-dipole interaction among the dipole-active vibrations of different molecules.

If $\boldsymbol{R}(t)$ is the intramolecular distance, and $\boldsymbol{d}$ and $\boldsymbol{d}^{\prime}$ are the operators of dipole moments of the molecules, the interaction Hamiltonian reads

$$
\begin{equation*}
\hat{V}=-|\boldsymbol{R}(t)|^{-5}\left[3(\boldsymbol{d} \cdot \boldsymbol{R}(t))\left(\boldsymbol{d}^{\prime} \cdot \boldsymbol{R}(t)\right)+\left(\boldsymbol{d} \cdot \boldsymbol{d}^{\prime}\right) R^{2}(t)\right] . \tag{2.21}
\end{equation*}
$$

The vibrational exchange is therefore mediated by the electromagnetic interaction, which occurs via emission and absorption of a photon near the point of closest approach, as shown in Fig. 2.7.

Let us assume that the main contribution to the vibrational energy exchange between identical colliding molecules comes from the dipole-dipole interaction of their identical modes that have the largest dipole moment. We may consider for example the $\nu_{3}$ modes of $S F_{6}$ molecules. Later on we show that if the typical interaction time is much shorter than the typical time of


Fig. 2.7. Vibrational energy exchange mediated by a photon. (a) Vibrational and rotational spectra of symmetric polyatomic molecules are mixed and form a common ro-vibronic spectrum. (b) Sketch of a transition accompanied by change of ro-vibronic states of colliding molecules. The dipole-dipole interaction is equivalent to a photon exchange. (c) The compound spectrum of the colliding molecules is a complex multilevel system.
the vibrational stochastization but much shorter than the inverse spacing among the neighboring quantum states (in units of $\hbar=1$ ), the equation for the ensemble average populations can be derived in the quasistatic assumption. In other words, one can find the parameters governing the population dynamics for the time-independent interaction and then allow for their slow time variation.

We mark the probability amplitudes by two subscripts: Latin indices denote the total energy of each molecule in terms of the number of quanta of the dipole-active mode, and Greek indices enumerate the sublevels of these levels emerging from the anharmonic splitting shown in Fig. 2.3.

The pair of colliding molecules comprise a compound system. At large distances when the molecules do not interact with each other, the state of the compound system $|n, m, \alpha, \beta\rangle=|n, \alpha\rangle \times|m, \beta\rangle$ is the direct product of the states of each molecule, and the amplitude of this state $\Psi_{n, m, \alpha, \beta}=\psi_{n, \alpha} \psi_{m, \beta}$ is the product of the amplitudes. At closer distances, for interacting molecules this is no longer the case, since the interaction (2.21) mixes up the states of different molecules and destroys the separability of their individual motions.

We therefore have to solve the Schrödinger equation for the entire compound system.


Fig. 2.8. Vibrational exchange in polyatomic molecules. (a) A photon exchange results in opposite vibrational transitions in colliding molecules. (b) Many photon exchange acts occur during the elementary act of a strong collision. (c) Transitions among the states of a compound system of two molecules is a process in a complex quantum system.

For interacting molecules we introduce new variables. By the number $N=(n+m) / 2$ we denote the conserving quantity - half of the total energy of the two molecules. Half of the energy difference $M=(n-m) / 2$ changes in the course of collision as well as the populations of the sublevels $\alpha$ and $\beta$ of the anharmonic splitting. In Fig. 2.8 we illustrate the correspondence between the transitions in each of the molecules and the transitions in the compound system. We denote by $\Delta_{N, M, \alpha, \beta}=E_{n, \alpha}-E_{m, \beta}$ the total energy of the non-interacting molecules and introduce the matrix elements of the interaction (2.21)

$$
\begin{equation*}
V_{N, M, \alpha, \beta}^{N, M+1, \alpha^{\prime}, \beta^{\prime}}=V_{n, m, \alpha, \beta}^{n+1, m-1, \alpha^{\prime}, \beta^{\prime}} ; \quad V_{N, M, \alpha, \beta}^{N, M-1, \alpha^{\prime}, \beta^{\prime}}=V_{n, m, \alpha, \beta}^{n-1, m+1, \alpha^{\prime}, \beta^{\prime}} \tag{2.22}
\end{equation*}
$$

In these notations the Schrödinger equation for the compound system reads

$$
\begin{align*}
i \hbar \dot{\Psi}_{N, M, \alpha, \beta}= & \Delta_{N, M, \alpha, \beta} \Psi_{N, M, \alpha, \beta}+\sum_{\alpha^{\prime}, \beta^{\prime}} V_{N, M, \alpha, \beta}^{N, M+1, \alpha^{\prime}, \beta^{\prime}} \Psi_{N, M+1, \alpha^{\prime}, \beta^{\prime}} \\
& +\sum_{\alpha^{\prime}, \beta^{\prime}} V_{N, M, \alpha, \beta}^{N, M-1, \alpha^{\prime}, \beta^{\prime}} \Psi_{N, M-1, \alpha^{\prime}, \beta^{\prime}} \tag{2.23}
\end{align*}
$$

The random matrix model is a natural selection for the interaction $\hat{V}$ in the chosen representation, since the anharmonic vibrational motion of each molecule is already chaotic, and the interaction would just add more complexity to the motion of the compound system.

### 2.2.2 Chemical Reactions

Chemical reactions of polyatomic molecules in the gas phase are a very important particular case of collisions. From the quantum mechanical point of view they are processes in multilevel quantum systems. We illustrate this in several examples.

## Harpoon Reaction

Consider an atom $A$ with a low ionization potential which is at a small distance $R$ from a diatomic molecule $B C$ which has a large electron affinity. The outer atomic electron can transfer to the molecule, which yields an ionmolecular compound $A^{+}+B C^{-}$. In Fig. 2.9 we illustrate this process and show the corresponding change $\Delta U\left(R_{B C}\right)$ of the internuclear potential energy $U\left(R_{B C}\right)$, which depends on the distance $R_{B C}$ between the atoms $B$ and $C$. The change of the potential results in a change of the molecular vibrations.

We express the wavefunction of the compound

$$
\begin{equation*}
\Psi=\sum_{n} \psi_{n}(t) \phi_{n} \Phi_{A}+\sum_{k} \psi_{k}(t) \phi_{k} \Phi_{B C} \tag{2.24}
\end{equation*}
$$

in terms of the electronic wavefunctions $\Phi_{A}\left(r, R, R_{B C}\right)$ localized at the atom and $\Phi_{B C}\left(r, R, R_{B C}\right)$ localized at the molecule, and the molecular wavefunctions $\phi_{n}\left(R_{A B}\right)$ and $\phi_{k}\left(R_{A B}\right)$ of the vibrational energy eigenstates $|n\rangle$ and $|k\rangle$ in the potentials $U\left(R_{B C}\right)$ and $U\left(R_{B C}\right)+\Delta U\left(R_{B C}\right)$ respectively. The coefficients $\psi_{n}(t)$ stand for the time-dependent probability amplitudes of the states when the electron is associated with the atom $A$ and the molecule $B C$ is at the $n$ vibrational level of the potential $U\left(R_{B C}\right)$, whereas the coefficients $\psi_{k}(t)$ correspond to the electron associated with the molecule and the vibrational state $k$ in the potential $U\left(R_{B C}\right)+\Delta U\left(R_{B C}\right)$.

We substitute the ansatz (2.24) into the Schrödinger equation and arrive at the equation for the amplitudes $\psi_{n}(t)$ and $\psi_{k}(t)$ :

$$
\begin{align*}
& i \hbar \dot{\psi}_{n}=E_{n} \psi_{n}+\sum_{k} V_{n, k} \psi_{k} \\
& i \hbar \dot{\psi}_{k}=E_{k} \psi_{k}+\sum_{m} V_{k, n} \psi_{n} \tag{2.25}
\end{align*}
$$



Fig. 2.9. The harpoon reaction. Atom $A$ loses an electron which is captured by the molecule $B C$ at a large distance $R$. The interatomic potential of the molecular ion differs by a value $\Delta U\left(R_{B C}\right)$ from the molecular potential $U\left(R_{B C}\right)$, and this difference serves as a perturbation for the vibrational motion of the ion.

Here $E_{n}$ and $E_{k}$ are the energies of the vibrational states $|n\rangle$ and $|k\rangle$. The probability amplitudes of the transitions are given by the overlap integral

$$
\begin{equation*}
V_{n, k}=\int \phi_{n}\left(R_{A B}\right) \phi_{k}\left(R_{A B}\right) \Phi_{B C}\left(\boldsymbol{r}, R, R_{B C}\right) \Phi_{A}\left(\boldsymbol{r}, R, R_{B C}\right) d^{3} r d R_{B C} \tag{2.26}
\end{equation*}
$$

where $\boldsymbol{r}$ denotes the coordinates of the transferred electron.
For a polyatomic molecule $M$ one has to replace the integration over the interatomic distance $R_{B C}$ by integration over all $N$-dimensional space of the internal degrees of freedom $d^{N} \boldsymbol{R}_{M}$. For the vibrational energies corresponding to the chaotic motions the overlap integrals

$$
\begin{equation*}
V_{n, k}=\int \Phi_{M}\left(\boldsymbol{r}, R, \boldsymbol{R}_{M}\right) \Phi_{A}\left(\boldsymbol{r}, R, \boldsymbol{R}_{M}\right) \phi_{n}\left(\boldsymbol{R}_{M}\right) \phi_{k}\left(\boldsymbol{R}_{M}\right) d^{N} \boldsymbol{R}_{M} d^{3} r \tag{2.27}
\end{equation*}
$$

are complex functions of the indices $n$ and $k$ and the matrix $V_{n, m}$ can therefore be replaced by a random matrix. We note that if the atom $A$ moves with respect to the molecule, the vibrational energy eigenstates and the transition amplitudes $V_{n, m}$ become time dependent.

## Sub-barrier Reaction

Another example concerns a chemical reaction which occurs in the framework of a single adiabatic electronic term, that is it involves only the vibrational motion of molecules. Consider an atom $A$ slowly approaching a rapidly oscillating molecule $B C$. The presence of the atom effects the frequency of the molecular vibrations, which gradually change as the atom closes with the molecule. After the atom has passed, the vibrational frequency returns to its initial value. This adiabatic change of frequency is a time-dependent perturbation, and in principle it is able to break the molecular bond $B-C$, although with a low probability. The energy conservation law implies that the breaking of the $B-C$ bond is followed by the formation of another bond, say $A-B$.


Fig. 2.10. Vibrational wavefunctions (a), potential surfaces (b), and trajectories (c) of a tunneling chemical reaction.

A consistent quantum description of the process should consider the process as a tunneling through the potential barrier along the reaction coordinate, which is shown in Fig. 2.10. Let us consider possible trajectories of the molecular motion in such a term. The trajectory starts at the potential valley that corresponds to the atom $A$ approaching the molecule $B C$. The oscillations of the trajectory represent the molecular vibrations. A non-reactive process corresponds to the trajectory 1 that leaves the region of interaction via the same potential valley, although probably with a different ratio between the vibrational and the translational energies. The chemical reaction corresponds to the trajectory 2 , which escapes the interaction region via the other potential valley. It also includes a tunneling part 3 . The oscillations of this trajectory now represent the vibrations of the molecule $A B$ that has been formed in the course of the chemical reaction.

We assume that the relative motion of the atom and the molecule is slow and write the Schrödinger equation

$$
\begin{align*}
& i \hbar \dot{\psi}_{n}=E_{n}(z) \psi_{n}+\sum_{k} V_{n, k}(z) \psi_{k} \\
& i \hbar \dot{\psi}_{k}=E_{k}(z) \psi_{k}+\sum_{m} V_{k, n}(z) \psi_{n} \tag{2.28}
\end{align*}
$$

for the amplitudes $\psi_{n}$ of the vibrational states $|n\rangle$ of the molecule $B C$, and the amplitudes $\psi_{k}$ of the vibrational states $|k\rangle$ of the reaction product $A B$. Here $z=z(t)$ is the slowly changing distance between the atom and the molecule, and $V_{n, k}$ represents the probability amplitudes of the tunneling through the potential barrier given by the overlap integrals of the corresponding wavefunctions (Fig. 2.10). The latter evidently depend on the distance $z$ between the atom and the molecule.

When we consider the reactive collisions of polyatomic molecules, the energies $E_{n}(z)$ and $E_{k}(z)$ acquire a complex structure as do the overlap integrals $V_{n, k}$. Such reactions occur in a multidimensional space of the configurations which is difficult to depict, although the main features of the process, namely the slow approach of the reagents and the tunneling along the reaction coordinate remain qualitatively the same. This tunneling may however occur in many different points of configuration space that correspond to different shapes of the molecules at the moment of the reaction. Moreover, the motion of polyatomic molecules is usually chaotic, and therefore the random matrix model is natural for the transition amplitudes $V_{n, k}$.

In order to calculate the reaction rate or the energy distribution of the products we first have to solve (2.28) and determine which combinations of the spectra $E_{n}(z) ; E_{k}(z)$ and the interaction $V_{n, k}$ govern the reaction properties.

### 2.2.3 Intermolecular Conversion and Photochemistry

From the conceptual point of view chemical reactions are very similar to the problems of intermolecular transfer of electronic energy to vibrational degrees of freedom, to photochemical processes, and to photon-assisted reactive collisions. We illustrate this with some examples.

## Intermolecular Conversion

The adiabatic separation of electronic and nuclear motion is valid only far from the points of term crossing. In the vicinities of these points the nuclear motion may cause electronic transitions, which implies energy exchange between electronic and vibrational degrees of freedom. This perturbation is of a higher order in the Born-Oppenheimer parameter $\Lambda_{B O}=\left(m_{e} / m_{p}\right)^{1 / 4}$ and normally is ignored, unless two or more electronic terms happen to be close
to each other at a certain position of the nuclei in the molecule. In the latter case the non-adiabatic interactions have to be taken into account.

Let us consider two crossing terms shown in Fig. 2.11. We take the diabatic potentials and construct the basis set from the electronic wavefunctions $\Phi_{1}(\boldsymbol{r}, \boldsymbol{R}), \Phi_{2}(\boldsymbol{r}, \boldsymbol{R})$ and the vibrational wavefunctions $\phi_{n}(\boldsymbol{R}), \phi_{k}(\boldsymbol{R})$ in the first and the second potentials respectively. Here $\boldsymbol{r}$ denotes the positions of


Fig. 2.11. Crossing of molecular terms (a). Diabatic terms (dotted line) correspond to a certain electronic state of a molecule and serve as potentials for the vibrational motion. Non-adiabatic interactions result in a reconstruction of the electronic terms and yield adiabatic terms (solid line) for which the electronic state is entangled with the nuclear position (sometimes also with the momentum). (b) Phenomena that occur near the crossing point can be considered in the diabatic vibrational basis as dynamical processes in a complex multilevel system.
all molecular electrons, and $\boldsymbol{R}$ stands for the position of all nuclei. For the corresponding amplitudes $\psi_{n}$ and $\psi_{k}$ we write the Schrödinger equation

$$
\begin{align*}
& i \hbar \dot{\psi}_{n}=E_{n} \psi_{n}+\sum_{k} V_{n, k} \psi_{k} \\
& i \hbar \dot{\psi}_{k}=E_{k} \psi_{k}+\sum_{n} V_{k, n} \psi_{n} \tag{2.29}
\end{align*}
$$

where $V_{n, k}$ now denotes the matrix element of the non-adiabatic interaction

$$
\begin{equation*}
V_{n, k}=-\int d^{p} r d^{N} \Phi_{1}(\boldsymbol{r}, \boldsymbol{R}) \phi_{n}(\boldsymbol{R})\left(\boldsymbol{\nabla}_{R} \Phi_{2}(\boldsymbol{r}, \boldsymbol{R}) \cdot \boldsymbol{\nabla}_{R} \phi_{k}(\boldsymbol{R})\right) . \tag{2.30}
\end{equation*}
$$

Here $p$ and $N$ denote the number of electronic and vibrational degrees of freedom respectively.

For polyatomic molecules the vibrational energies $E_{n}$ and $E_{k}$ and the matrix elements $V_{n, k}$ can be modeled, as earlier, by random values.

## Photochemistry and Laser-Induced Term Crossing

Optical or ultraviolet laser fields induce transitions among the electronic terms of molecules that usually result in chemical transformations. Strong fields may also affect the vibrational dynamics. In the last case one can describe the system by (2.29), where instead of the matrix elements of nonadiabatic interaction (2.30) we have to substitute the probability amplitudes of laser induced transitions

$$
\begin{equation*}
V_{n, k}=-\mathcal{E} e \int d^{p} r d^{N} \Phi_{1}(\boldsymbol{r}, \boldsymbol{R}) \phi_{n}(\boldsymbol{R}) \boldsymbol{r} \Phi_{2}(\boldsymbol{r}, \boldsymbol{R}) \phi_{k}(\boldsymbol{R}) \tag{2.31}
\end{equation*}
$$

Here $\mathcal{E}$ is the electric field amplitude, and $e$ is the charge of the electron.
This analogy becomes evident, when we consider a compound system of the molecule and the quantized laser field. Indeed, at some position of the nuclei, the energy of the electron cloud in the ground electronic term equals the energy of the excited electronic state minus the energy of one photon. One can say that at this point, the two potential curves of the compound system cross each other. Coupling occurs due to the dipole interaction with the laser field, which is taken into account by the matrix elements (2.31).

If the typical size of the interaction exceeds the typical energy spacing among the neighboring vibrational states, this process cannot be considered perturbativly. For diatomic molecules it requires a very high laser intensity, at which the typical interaction $\mathcal{E} d$ is of the order of magnitude or larger than the vibrational quantum. However this requirement is not as restricting for polyatomic molecules, which have a much higher density of vibrational states: even moderately intense laser fields may result in the coupling of electronic and vibrational motions. The complex structure of the vibrational wavefunctions of the excited polyatomic molecules results in an irregular dependence of $V_{n, k}$ on $n$ and on $k$, and therefore justifies the utilization of a random matrix model.

In some cases more than two electronic terms may get involved in the dynamical process at moderate laser intensities. Then for the description of such a system one has to employ (2.16) where $I$ denotes the number of absorbed quanta, and $m$ enumerates the vibrational eigenstates in each of the electronic terms involved.

Accommodation of the energy of an absorbed optical quantum by molecular vibrations is also an example of processes in multilevel systems. Consider a molecule initially in the ground vibrational state of the lowest electronic term which absorbs a quantum of radiation and makes a transition to an excited vibrational state of the higher electronic term. Thence, due to the non-adiabatic interaction it gets transferred back to the ground electronic term, although in a different vibrational state. This state may also be predissociating. The Schrödinger equation for the amplitudes $\psi_{k}$ of the vibrational states in the diabatic excited electronic term, the amplitudes $\psi_{n}$ of the highly excited vibrational states of the ground diabatic term, and for the amplitude
$\psi_{0}$ of the ground vibrational state read

$$
\begin{align*}
& i \hbar \dot{\psi}_{0}=E_{0} \psi_{0}+\sum_{k} \mathcal{E} d_{0, k} \psi_{k} \\
& i \hbar \dot{\psi}_{n}=E_{n} \psi_{n}+\sum_{k} V_{n, k} \psi_{k} \\
& i \hbar \dot{\psi}_{k}=\left(E_{k}-\hbar \omega\right) \psi_{k}+\sum_{n} V_{k, n} \psi_{n}+\mathcal{E} d_{k, 0} \psi_{0} \tag{2.32}
\end{align*}
$$

Here $V_{m, n}$ represents the matrix elements of the non-adiabatic transitions (2.30), and $\mathcal{E} d_{0, k}$ are the matrix elements (2.31) of the interaction with the electromagnetic field.

### 2.3 Rydberg Molecules

Multilevel quantum problems arise in a natural way when one considers highly electronically excited or so-called Rydberg molecules, that is molecules with one electron excited into a high Rydberg state. This electron possesses an energy close to the ionization potential and, in the language of classical mechanics, it moves along a long Kepler orbit in the electrostatic field of the molecular ionic core. Traditional adiabatic separation of the electronic and vibrational motion in molecules implies a large difference in energies, velocities and frequencies of the electrons and nuclei of molecules. This is not completely the case for the Rydberg molecules: the Kepler period might be comparable and even longer than the vibrational or rotational period of the ionic core, although the velocity of the Rydberg electron near the core is still much higher than the velocity of the nuclei. We therefore meet the situation which in a sense is inverse to the regular adiabatic case: the electron is a slowly moving subsystem, while the rotations and vibrations of the ionic core are the fast motion. We can therefore speak about the Rydberg sublevels $E_{n}$ of each ro-vibrational state $|k\rangle$ of the molecular ion. We give two examples of such molecules.

### 2.3.1 Subthreshold Photoionization

Consider a molecule in an intense laser field. The outer electron of the molecule makes a multiphoton transition to a highly excited Rydberg state and starts its motion along a Kepler orbit. Returning back to the vicinity of the molecular ion it encounters a potential that differs from the Coulombic case, and therefore experiences scattering at the asymmetric part of the potential. The scattering may be either elastic or non-elastic. The non-elastic scattering is associated with a transfer of energy and angular momentum from the electron to the molecular ion, and hence the electron changes its


Fig. 2.12. Multiphoton ionization of a molecule via Rydberg states.

Rydberg orbit (see Fig. 2.12), whereas the molecular ion finds itself in another rotational or vibrational state. The scattered electron with an energy less than the ionization potential returns back again after a long Kepler period. During this time the molecular ion in the center vibrates and rotates, and there is sufficient time for the position of the nuclei to change. Therefore the returning electron experiences quite a different scatterer as compared to the initial scattering event. Thus the sequence of scattering events is a kind of random process which has a direct analogy to the quantum Sinai billiard.

For the basis set of the problem we employ the products $|n\rangle|k\rangle$ of the Rydberg states $|n\rangle$ of the hydrogen atom and the rotational (and/or vibrational) states $|k\rangle$ of the molecular ion, and write down the Schrödinger equation

$$
\begin{align*}
i \hbar \dot{\psi}_{0} & =E_{0} \psi_{0}+\sum_{k, n} \tilde{V}_{0}^{k, n} \psi_{k, n} \\
i \hbar \dot{\psi}_{n, k} & =\left(E_{n}+E_{k}\right) \psi_{k, n}+\sum_{k^{\prime}, n^{\prime}} V_{n, k}^{k^{\prime}, n^{\prime}} \psi_{k^{\prime}, n^{\prime}}+\tilde{V}_{k, n}^{0} \psi_{0} \tag{2.33}
\end{align*}
$$

for the amplitudes $\psi_{k, n}$ of these states. Here $\tilde{V}$ stands for the amplitude (2.18) of multiphoton transitions from the ground state $|0\rangle$. The matrix elements $V_{k, n}^{k^{\prime}, n^{\prime}}$ allow for the probability amplitudes of scattering from the state $\left|k^{\prime}\right\rangle\left|n^{\prime}\right\rangle$ to the state $|k\rangle|n\rangle$.

We note one interesting peculiarity of this process. The Kepler period may change considerably with the change in the energy of the Rydberg electron: the closer the energy to the ionization limit the longer the period. In such a situation it is natural to consider the evolution of the system not as a function of time, but as a function of the number of scattering events. One determines the corresponding transformation by multiplying the scattering amplitudes $V_{k, n}^{k^{\prime}, n^{\prime}}$ in (2.33) by the phase factor $\exp \left[-i\left(E_{k^{\prime}}+E_{n^{\prime}}\right) T\left(n^{\prime}\right)\right]$, where $T\left(k^{\prime}\right)$ is the Kepler period corresponding to the energy $E_{k^{\prime}}$. This phase factor changes drastically from one state to another, and it makes the matrix elements of the neighboring states very different in phase and thereby justifies the use of the random matrix model for this process.

We also note one important application of this process: by measuring the yield of zero kinetic energy electrons as a function of the laser frequency one obtains the high-resolution rotational and vibrational spectra of molecular ions.

### 2.3.2 Collisional Ionization

In Fig. 2.13 we show another example of a non-adiabatic process: the transfer of vibrational energy to the Rydberg electron in the course of associative ionization. An excited noble gas atom interacts with a metal atom. At interatomic distances shorter than the orbit size of the excited electron the metal atom sees only the ionic core of the noble atom, since the electron most of the time stays at a longer distance. From the chemical point of view the core resembles a halogen atom, that is an element with one vacancy in the outer electronic shell. Therefore the reaction between the metal and the noble gas core is of a harpoon type: the outer electron of the metal jumps to the core and fills the vacancy in the outer shell of the noble gas ion. This is essentially the formation of an ionic bond, with the two atoms comprising an ionic molecular core, which generally speaking is in a vibrationally excited state.

At the same time the excited electron of the noble gas atom feels a displacement of the charge center toward the metallic atom. If the distance between the electron and the charge center increases, the electron finds itself in a high Rydberg orbit of the molecular ion. The Kepler trajectory of the electron goes through the space occupied by the molecular core, and therefore the Rydberg electron experiences a sequence of scattering analagous to the case considered just before. In the course of non-elastic collisions the vibrational energy of the molecular core can be transferred to the Rydberg electron, which may result in ionization. The Schrödinger equation describing the process

$$
\begin{equation*}
i \hbar \dot{\psi}_{n, k}=\left(E_{n}+E_{k}\right) \psi_{k, n}+\sum_{k^{\prime}, n^{\prime}} V_{n, k}^{k^{\prime}, n^{\prime}} \psi_{k^{\prime}, n^{\prime}} \tag{2.34}
\end{equation*}
$$

coincides with (2.33) when the laser field is at zero strength, $\tilde{V}=0$.


Fig. 2.13. Associative ionization. The outer electron of a metal atom experences the vacancy in the outer shell of the noble gas core and forms a chemical bond. The excited electron of the noble gas atom finds itself in a high Rydberg orbit of the molecular ion. The vibrational energy of the molecular core is transferred to the Rydberg electron in the course of collisions between the electron and the core. The ionization therefore starts not at the incedent energy $E_{k}$ equal to the adiabatic threshold, but at a lower (subthreshold) energy.

### 2.4 Atomic and Molecular Clusters

Clusters are intermediate objects between molecules and solids, and the physics of clusters inherits complexity from both molecular and solid state physics. Indeed, the typical distance among the electronic states of a $N$ atom cluster is $N$ times smaller than for a single atom, and hence for $N \sim \Lambda_{B O}^{-1}=\left(m_{n} / m_{e}\right)^{1 / 4} \sim 10$ the adiabatic separation of the nuclear and electronic motion typical of molecules fails. In other words, for numbers of atoms comparable to or higher than the inverse of the Born-Oppenheimer parameter the nonadiabatic corrections are of the order of the typical distance between neighboring electronic states, and hence the electrons no longer follow the nuclear motion.

The concept of quasiparticles is the key idea of solid state physics. The strong interaction of electrons and nuclei comprising solids can be treated in
terms of almost free elementary excitations such as phonons, excitons, holes and electrons with an effective mass, etc. However, for clusters of sizes less than a typical free path of the quasi particles $(N \sim 1000)$ this concept does not yield any considerable simplification of the many-body problem: quantum interference effects resulting from the confinement, and strong interaction on the cluster boundaries make the quasiparticle approach as difficult as the quantum chemical one.

## Ensemble of interacting atoms

a)

excitation
exchange


Cluster
b)

$n=5$

$n=1$

Fig. 2.14. Reconstruction of the energy spectrum of an ensemble of two-level systems as a result of the excitation exchange interactions. (a) Energy spectrum of the ensemble of atoms. The mean interatomic distance is large and the spectrum is almost degenerate. (b) Energy spectrum of a cluster formed as a result of the reconstruction of the atomic spectrum due to strong interatomic interaction.

Modeling of interactions by random matrices appears in a natural way in the physics of atomic and molecular clusters. Indeed if, as the result of mutual attraction, a group of atoms or molecules is stuck together, the spectrum of the energy eigenstates of such aggregation changes considerably as compared to the initial spectrum of the non-interacting particles. We illustrate this in a simple example presented in Fig .2 .14 by showing a reconstruction of the spectrum of an ensemble of two-level systems resulting from the excitationexchange interaction among them. It might be the dipole-dipole interaction, for instance.

At finite temperatures various configurations of the cluster are possible, and each of the geometrical arrangements of the atoms corresponds to different interatomic distances, and therefore to different sizes of the pair interactions. Let us assume that only one atom was originally excited. Then, as a result of the interaction the excitation can pass to other atoms, and hence the energy eigenstates of the cluster do not correspond to a certain location of the excitation. If nevertheless for the basis set of the problem we chose the states corresponding to all possible locations of the excitation at a single atom, the interaction Hamiltonian becomes a matrix with the elements depending on the particular configuration of the cluster. For a random arrangement of atoms this matrix is random.

In spite of the fact that the interaction modifies considerably the energy eigenstates of the system, the unperturbed Hamiltonian $\hat{H}_{0}$ still plays an important role. We therefore arrive at the necessity to describe the energy spectrum of a system perturbed by a random matrix $\hat{V}$. Which characteristics have to be employed for the purpose? The density of states $\pi^{-1} \operatorname{Im} \operatorname{Tr}(E-\hat{H}-$ $i 0)^{-1}$ is one of the most represenative quantities for a spectrum. It suggests the ensemble averaged density of states

$$
\begin{equation*}
g(E)=\left\langle\frac{1}{\pi} \operatorname{Im} \operatorname{Tr} \frac{1}{E-\hat{H}_{0}-\hat{V}-i 0}\right\rangle \tag{2.35}
\end{equation*}
$$

as one of the key characteristics of a complex system. Indeed, the state density is an ingredient of many experimentally observable properties such as optical absorption, scattering cross-section, etc.

### 2.4.1 Ground Electronic State of Hot Metallic Clusters

We dwell on the example of metallic clusters which are an important particular case of aggregations. Metallic clusters have common orbitals for the valence electrons. For minute clusters the shell structure of the orbitals effects strongly their basic physical properties: the Fermi energy, work function, electron affinity, and cohesion energy all change abruptly and considerably with an increase of the cluster size. At finite temperatures the thermal motion of atoms perturbs appreciably the shell structure, and thereby it modifies significantly properties of the clusters.

Let us consider this in the framework of the simplest one-particle approximation. The valence electrons of a cluster move in the self-consistent effective potential, which implies that each of the electrons feels an average field of the ionic core and all other electrons. The simplest model of this field is just a rectangular potential well shown in Fig. 2.15. For spherically symmetric clusters the angular momentum $\boldsymbol{l}$ of each of the electrons is conserved, and their angular motion can be taken into account by introducing the centrifugal force potential $U(r)=(\boldsymbol{l})^{2} / 2 m_{e} r^{2}$, acting on the motion in the radial direction $r$.


Fig. 2.15. Spherically symmetric potential of cluster electrons. (a) Rectangular well model of the work function and self-consistent local density potential. (b) Potential of the centrifugal forces. (c) Effective potential and a wavefunction. (d) Classical trajectory of an electron.

By solving the corresponding Schrödinger equation one finds the energy eigenfunctions

$$
\begin{equation*}
\psi_{n, l, m}=\frac{J_{l+1 / 2}\left(k_{n, l} r\right)}{\sqrt{r}} Y_{l, m}(\theta, \phi) \tag{2.36}
\end{equation*}
$$

in terms of the Bessel functions $J_{\nu}(x)$, and the spherical harmonics $Y_{l, m}(\theta, \phi)$. For deep potential wells with the help of the boundary condition at the outer border of the cluster

$$
\begin{equation*}
\left.\psi_{n, l, m}\right|_{r=R}=0 \tag{2.37}
\end{equation*}
$$

one determines the radial wavevectors $k_{n, l}$ and the corresponding energies

$$
\begin{equation*}
E_{n, l}=\frac{\hbar^{2} z_{n, l}^{2}}{2 m_{e} r_{S}^{2} N^{2 / 3}} \tag{2.38}
\end{equation*}
$$

provided the number $N$ of atoms in the cluster and the typical size $r_{S}$ of the atoms (Wigner-Seitz radius) are known. Here $z_{n, l}$ stands for the $n$-th root of the Bessel function $J_{l+1 / 2}\left(z_{n, l}\right)=0$. The states are degenerate with respect to the magnetic quantum number $m$.

More accurate models include a self-consistent consideration of the electronic states based on the local density functional approach. This implies
iterative numerical calculation of the electronic wavefunctions and the corresponding effective potential determined as a functional of the electron density distribution. This method results in more accurate values of the cluster electron energies (2.38).

At finite temperatures the thermal motion of atoms comprising the cluster destroys the spherical symmetry and raises the spherical degeneracy of the electron energy eigenstates. Since the atomic motion is much slower than the electron motion we can consider this perturbation as static. We note, however, that this assumption does not mean that the momenta of the electrons play no role in the perturbation. On the contrary, the nuclear momenta bring an important contribution via the non-adiabatic coupling (2.30) of the electronic states. The assumption only means that the time variation of the coupling is slower than its typical size divided by Planck's constant.


Fig. 2.16. The classical trajectory of an electron (dashed line) in a hot cluster. Scattering of the electron occures at impurities (i), at phonons (p), and at the perturbed surface (s) of the spherical "in average" cluster.

The random character of the thermal motion suggests the consideration of the perturbation as random. Indeed, in the classical limit the trajectory of an electron moving in the potential well resembles the trajectory of a particle in Sinai billiard theory (see Fig. 2.16): elastic scattering at thermal phonons at irregularities of the surface and at atomic cores makes the electronic motion chaotic. Therefore the corresponding quantum model relies on the assumption that the thermal motion results in a coupling of the electronic states (2.36) of the cluster via a random matrix with the mean squared matrix element $\left\langle V^{2}\right\rangle$ depending on temperature $T$. One can estimate this quantity by noticing that the electron-phonon coupling is the main mechanism of the temperaturedependent conductivity of metals, and the relaxation time of the electronic momentum $\tau(T)$ equals the inverse of the relaxation probability $2 \pi\left\langle V^{2}\right\rangle g_{f}$,
where $g_{f}$ is a typical one-electron state density near the Fermi energy of the cluster.

Many of the properties of clusters in the ground electronic state can be found with the help of the ensemble averaged density of the one-electron states (2.35) where the energies (2.38) give the unperturbed Hamiltonian $\hat{H}_{0}$. In particular, one finds the Fermi energy $\epsilon_{f}$ of the $N$-electron cluster from the normalization condition

$$
\begin{equation*}
N=\int_{-\infty}^{\epsilon_{f}\left(N,\left\langle V^{2}\right\rangle\right)} g(E) d E \tag{2.39}
\end{equation*}
$$

the total energy of the electronic cloud

$$
\begin{equation*}
\mathcal{E}\left(N,\left\langle V^{2}\right\rangle\right)=\int_{-\infty}^{\epsilon_{f}\left(N,\left\langle V^{2}\right\rangle\right)} E g(E) d E \tag{2.40}
\end{equation*}
$$

and the dissociation energies

$$
\begin{equation*}
D\left(N,\left\langle V^{2}\right\rangle\right)=E_{I}+\mathcal{E}\left(N-1,\left\langle V^{2}\right\rangle\right)-\mathcal{E}\left(N,\left\langle V^{2}\right\rangle\right) \tag{2.41}
\end{equation*}
$$

of an ensemble of clusters. Here $E_{I}$ is the ionization potential of atoms comprising the cluster. Therefore the establishing of a simple recipe for the calculation of the ensemble average density of states $g(E)$ for an arbitrary Hamiltonian $H_{0}$ is of practical importance.

### 2.4.2 Optical Properties of Clusters

The random matrix model can also be applied to the excited states of metallic clusters. One constructs these states by considering elementary excitations, that is by transferring an electron from one of the occupied one-electron states $\left|n_{e}, l_{e}, m_{e}\right\rangle$ to one of the empty states $\left|n_{h}, l_{h}, m_{h}\right\rangle$ above the Fermi energy. In other words the elementary excitation corresponds to an electron in the state $\left|n_{e}, l_{e}, m_{e}\right\rangle$ and a hole in the state $\left|n_{h}, l_{h}, m_{h}\right\rangle$. For the non-interacting electrons it yields the energies

$$
\begin{equation*}
E\left(n_{e}, l_{e}, m_{e}, n_{h}, l_{h}, m_{h}\right)=E_{n_{e}, l_{e}, m_{e}}-E_{n_{h}, l_{h}, m_{h}} \tag{2.42}
\end{equation*}
$$

of the excited states of clusters.
Electrons however interact with each other. Some part of this interaction has already been taken into account when the self-consistent effective potential has been employed as a model for the ground electron term. The remaining part $\hat{V}$ should allow for the deviation of the real interaction from the effective one. The main contribution $\hat{V}^{(1)}$ to this remainder originates from the coupling of the electron-hole pairs with the electromagnetic field. The excited electron is annihilated with the hole and creates an intermediate virtual photon. The cluster therefore returns to the ground electronic


Fig. 2.17. Reconstructed spectrum of the exited states of $N a_{21}^{+}$.
term. The photon, in its turn, creates another electron-hole pair by transporting an electron from another occupied one-electron level to a different empty state. Therefore all excited states $|S\rangle=\left|n_{e}, l_{e}, m_{e}, n_{h}, l_{h}, m_{h}\right\rangle$ interact among themselves via the ground electronic term and different states of the electromagnetic field. It yields a specific (factorized) structure of the interaction Hamiltonian: each matrix element $V_{S, S^{\prime}}^{(1)}=T_{S} \times T_{S^{\prime}}$ is a product of a quantity related to the starting state $\left|S^{\prime}\right\rangle$ and the same quantity related to the final state $|S\rangle$. One has to include this interaction in the Hamiltonian of the cluster, which yields a considerable modification of the energies (2.42). An example of such a reconstruction of a spectrum is shown in Fig. 2.17.

The rest $V_{S, S^{\prime}}^{(2)}$ of the remaining $V_{S, S^{\prime}}$ has a complicated structure. It originates from the scattering of the excited electron at the hole. The dynamics of the electron resemble the motion in a Sinai billiard with one difference: the scatterer, that is the hole, does not stay at the same position but also moves along a chaotic trajectory as the result of scattering by the electron. Therefore the matrix elements $V_{S, S^{\prime}}^{(2)}$ do not have a factorized structure. On the contrary, they depend in a complex and irregular way on the indices $S$ and $S^{\prime}$ and often are ignored. This is known as the random phase approximation. However, this interaction can also be taken into account within the random matrix model where the energies (2.42) and the interaction $V^{(1)}$ compose the unperturbed Hamiltonian, and the interaction $V^{(2)}$ gives a random perturbation. One can also include in the latter part the random coupling originating from the scattering of the electron and the hole at the atomic core phonons and irregularities of the surface by analogy to the coupling in the ground electronic term discussed in the previous section.

One ususally finds cross-sections of the linear absorption from the ground state $|0\rangle$ of a quantum system with Hamiltonian $\hat{H}$ with the help of the expression

$$
\begin{equation*}
\sigma(\hbar \omega)=\frac{1}{\pi} \operatorname{Im}\langle 0| \hat{d} \frac{1}{\hbar \omega-\hat{H}-i 0} \hat{d}|0\rangle, \tag{2.43}
\end{equation*}
$$

where $\hat{d}$ is the dipole moment operator along the electric field vector. The real part of the quantum average allows for the optical dispersion. For ensembles of clusters (2.43) yields

$$
\begin{equation*}
\left.\sigma(\hbar \omega)=\frac{1}{\pi} \operatorname{Im}\left\langle\left.\langle 0| \hat{d} \frac{1}{\hbar \omega-\hat{H}_{0}-\hat{V}-i 0} \hat{d} \right\rvert\, 0\right\rangle\right\rangle \tag{2.44}
\end{equation*}
$$

where $\hat{H}_{0}$ is the ensemble average Hamiltonian, and $\hat{V}$ stands for the random perturbation. More sophisticated combinations can be written for multiphoton cross-sections.

### 2.5 Some Other Examples

We consider here several examples from branches of physics other than atomic and molecular physics, chemistry, or cluster physics. These examples relate to fields with their own conceptual systems and corresponding well-developed mathematical techniques. However, it seems useful to trace the analogies among these problems in terms of interacting quantum states in order to gain a deeper insight into the common properties of complex quantum objects.

### 2.5.1 Ion Traps

Systems such as ions trapped in an oscillating electromagnetic field where classical and quantum effects are of the same order of magnitude represent an interesting branch of contemporary physics. From electrostatics we know that a given number of charged particles can never form an equilibrium configuration in the static electric field. On the other hand, in a time-dependent field a quasi-equilibrium and dynamically stable configurations may exist. Indeed, rapid oscillation of the field results in an effective Kapiza-Dirac potential, which may confine the motion of charged particles in a certain region of space. A quasi-equilibrium distribution arises as the result of the interplay between Coulomb repulsion of the ions and the confining action of the effective potential: the particles rapidly oscillate around certain points that either move slowly or are static. This dynamic equilibrium may however be destroyed by the influence of the small terms oscillating at double the driving field frequency that result from the nonlinearity of the system. These terms are usually ignored in the first-order approximation. One can take them into account with the help of the random matrix model by the analogy to polyatomic molecules in an external laser field.

Let us illustrate this with the example of $N$ ions in a time-dependent electric field with the potential

$$
\begin{equation*}
U(\boldsymbol{r}, t)=\alpha\left(x^{2}+z^{2}-2 y^{2}\right)+A\left(x^{2} \cos \Omega t-A y^{2} \cos \Omega t\right) \tag{2.45}
\end{equation*}
$$

where $\boldsymbol{r} \equiv(x, y, z)$ is the radius vector and $\Omega$ is the oscillation frequency of the electric field vector in the plane $(x, y)$. One sees that both static and time-dependent components of the potential satisfy the Laplace equation $\Delta U=0$, which means that the chosen potential exists. We start our analysis in terms of classical mechanics and quantize the motion at a later stage, since the result of such consideration does not differ much from the result of a consistent quantum approach.


Fig. 2.18. Ions in the effective potential of an RF trap. (a) A static electric field (shown by the field lines) rotates with an angular velocity $\Omega$ around the $z$-axis. (b) The effective potential resulting from the rotation results in confinement of atomic ions, which therefore form a crystal. Higher order harmonics of the effective potential play the role of an external periodic force. (c) Such a perturbation provokes transitions among the eigenenergy quantum states in the effective potential.

The classical Hamiltonian of $N$ ions in the field reads

$$
\begin{equation*}
H\left(\left\{\boldsymbol{p}_{n}\right\},\left\{\boldsymbol{r}_{n}\right\}\right)=\sum_{n=1}^{N}\left(\frac{\boldsymbol{p}_{n}^{2}}{2 m}+e U\left(\boldsymbol{r}_{n}, t\right)\right)+\sum_{n \neq m} \frac{e^{2}}{\left|\boldsymbol{r}_{n}-\boldsymbol{r}_{m}\right|}, \tag{2.46}
\end{equation*}
$$

where $e, \boldsymbol{p}_{n}$ are the ion charges and classical momenta respectively, and $m$ is the ion mass.

In a high-frequency field the motion of each particle consists of two components

$$
\begin{equation*}
\boldsymbol{r}_{n}=\boldsymbol{\rho}_{n}(t)+\boldsymbol{\xi}_{n}(t), \tag{2.47}
\end{equation*}
$$

where the $\boldsymbol{\xi}_{n}(t)$ rapidly oscillates at the frequency of the field, and $\boldsymbol{\rho}_{n}(t)$ is a slow component. The typical value $\xi_{t}$ of the rapid component is usually much smaller than the typical value of the slow component $\rho_{t}$, and therefore the ratio $\xi_{t} / \rho_{t}$ is a small parameter. This is not the case for the velocities

$$
\begin{equation*}
\dot{\boldsymbol{r}}_{n}=\dot{\boldsymbol{\rho}}_{n}(t)+\dot{\boldsymbol{\xi}}_{n}(t), \tag{2.48}
\end{equation*}
$$

since $\xi$ rapidly changes with time. In fact, we have $\dot{\boldsymbol{\rho}}_{n}(t) \sim \dot{\boldsymbol{\xi}}_{n}(t)$.
We take into account that $\boldsymbol{p}_{n}=m \dot{\boldsymbol{r}}_{n}=m \dot{\boldsymbol{\rho}}_{n}(t)+m \boldsymbol{\xi}_{n} \equiv \mathcal{P}_{n}+\boldsymbol{\pi}_{n}$ and substitute (2.47), (2.48) into (2.46), which yields

$$
\begin{align*}
H\left(\left\{\boldsymbol{p}_{n}\right\},\left\{\boldsymbol{r}_{n}\right\}\right)= & \sum_{n=1}^{N}\left(\frac{\boldsymbol{\mathcal { P }}_{n}^{2}+\boldsymbol{\pi}_{n}^{2}}{2 m}+\frac{\left(\mathcal{P}_{n} \cdot \boldsymbol{\pi}_{n}\right)}{m}+e U\left(\boldsymbol{\rho}_{n}+\boldsymbol{\xi}_{n}, t\right)\right) \\
& +\sum_{n \neq m} \frac{e^{2}}{\left|\boldsymbol{\rho}_{n}+\boldsymbol{\xi}_{n}-\boldsymbol{\rho}_{m}+\boldsymbol{\xi}_{m}\right|} . \tag{2.49}
\end{align*}
$$

The rapidly oscillating components of the Hamiltonian (2.49) allow for the rapid motion, and the slowly varying components for the slow motion. We note that the term $\pi_{n}^{2}$ contains a slow component, namely the average kinetic energy of the rapid oscillations. This component plays the role of the potential energy for the slow motion, and therefore together with the component of the potential energy depending on the slow coordinates it comprises the Kapiza-Dirac potential. Indeed, the time derivatives of the rapid momenta, say in the $x$-direction, equal the rapid forces $\boldsymbol{\pi}_{n}^{(x)}=2 A x \cos \Omega t$, and hence $\boldsymbol{\pi}_{n}^{(x)}=2 A x / \Omega^{-1} \sin \Omega t$, which yields the contribution $\left\langle\left(\boldsymbol{\pi}_{n}^{(x)}\right)^{2} / 2 m\right\rangle_{t}=$ $2(A x / \Omega)^{2} / m\left\langle\sin ^{2} \Omega t\right\rangle_{t}=(A x)^{2} / m$ to the effective potential.

The Hamiltonian for the slow motion reads

$$
\begin{align*}
H\left(\left\{\boldsymbol{p}_{n}\right\},\left\{\boldsymbol{r}_{n}\right\}\right)= & \sum_{n=1}^{N}\left(\frac{\mathcal{P}_{n}^{2}}{2 m}+e \alpha\left(x_{n}^{2}+z_{n}^{2}-2 y_{n}^{2}\right)+\frac{e A^{2}}{m \omega^{2}}\left(x_{n}^{2}+y_{n}^{2}\right)\right) \\
& +\sum_{n \neq m} \frac{e^{2}}{\left|\boldsymbol{\rho}_{n}-\boldsymbol{\rho}_{m}\right|}, \tag{2.50}
\end{align*}
$$

and hence for $2 \alpha<A^{2} / m \omega^{2} m$, the effective potential becomes binding as shown in Fig. 2.18 Therefore the slow motion is confined and from the viewpoint of this motion the system possesses an equilibrium configuration. One
can say that the time-dependent electrical field creates a parabolic trap for the ions. Quantization of the slow motion in the trap yields a set of slowenergy eigenstates.

However, these states are not real eigenstates of the system, since the higher-order expansion terms

$$
\begin{equation*}
V\left(\left\{\rho_{n}\right\}\right)=\sum_{k, l, n \neq m} \frac{1}{l!k!} \frac{\partial^{l+k}}{\partial^{l} \boldsymbol{\rho}_{n} \partial^{k} \boldsymbol{\rho}_{m}} \frac{e^{2}}{\left|\boldsymbol{\rho}_{n}-\boldsymbol{\rho}_{m}\right|} \boldsymbol{\xi}_{m}^{l} \boldsymbol{\xi}_{n}^{k} \tag{2.51}
\end{equation*}
$$

of the Coulomb interaction in $\boldsymbol{\xi}=\boldsymbol{\rho}\left(A / m \Omega^{2}\right) \sin \Omega t$ result in a timedependent perturbation and therefore they induce transitions among the eigenstates of slow energy by analogy to multiphoton transitions in atoms or molecules. The frequencies of the perturbation are multiples of the driving frequency $\Omega$. We note that the consistent quantum consideration yields an expansion slightly different from (2.50); however this difference should not play a significant role for relatively heavy ions and the relatively large typical geometrical size $\alpha^{-1 / 2}$ of the potential.

Classical motion of the ions in the trap is multidimensional and nonlinear. Therefore at a certain level of the excitation it becomes chaotic. In this case one can consider the matrix elements of perturbation (2.50) as random, and write the Schrödinger equation (2.6) by analogy to polyatomic molecules where one also should allow for the transitions with the energy changing by $2 \hbar \Omega, 3 \hbar \Omega$, etc.

### 2.5.2 Disordered Solids and Surfaces

The description of quantum effects in terms of random matrices is a powerful method to tackle a number of problems in solid state physics and surface science. This approach arises in a natural way if one considers the properties of disordered systems like alloys, glasses, non-ideal metals, or systems with chaotic billiard-like dynamics such as quantum dots. It turns out that for such objects interference phenomena play a crucial role by affecting considerably their macroscopic properties. The basic physical processes take place on a limited size-scale which exceeds the atomic sizes, but remains much smaller than a typical macroscopic scale. Therefore these phenomena are called "mesoscopic", and comprise a special field of solid state physics.

Mesoscopic phenomena in solids have as a rule two peculiarities. They usually occur in a coordinate space of three or less dimensions and remain translationally invariant "on average". Therefore the properties of large species result from a kind of superposition of the mesoscopic properties of many statistically equivalent parts of smaller sizes, attached to each other consecutively. It allows one to take advantage of special basis sets, like plane waves for example, to use the scaling hypothesis, or to factorize the propagators of the particles.

## Disordered Chains

We illustrate the typical formulation of such a problem by the classical example of the model which considers a single particle in a disordered potential. We take the simplest one-dimensional version of the potential $U(x)=\sum_{n} k_{n} \delta(x-n a)$, which is a sequence of equidistant Dirac $\delta$-functions shown in Fig. 2.19 of different "depth" $k_{n}$, where the $k_{n}$ values are random. The particle may be located at a given knot $n$ of the potential or tunnel at a distance $a$ to a neighboring knot with a transition amplitude $V$. The Schrödinger equation for the probability amplitudes $\psi_{n}$ of the states $|n\rangle$ corresponding to the particle at a given knot $n$ can be written in the form

$$
\begin{equation*}
E \psi_{n}=\delta k_{n} \psi_{n}+V \psi_{n+1}+V \psi_{n-1} \tag{2.52}
\end{equation*}
$$

where $E$ and $\delta k_{n}$ are the deviations of an eigenenergy and the individual depth from the mean depth respectively, and $V$ is the tunneling amplitude.
a)

b)


Fig. 2.19. A random sequence (a) of $\delta$-potentials (Lloyd chain) and a more complicated potential (b) with microstructure.

The most remarkable property of (2.52) is the so-called Anderson localization phenomenon. Any energy eigenstate $|E\rangle$ of this equation is located mainly at a limited number of states $|n\rangle$, which means that the coefficients $c_{E}^{n}$ of the expansion $|E\rangle=\sum c_{E}^{n}|n\rangle$ drop rapidly when $n$ deviates from the position $n_{\max }(E)$ of the knot corresponding to the maximum absolute value of the amplitude $\psi_{n}$. The dependence of $c_{E}^{n}$ on $n$ is not monotonic and on the contrary is rather random. However, the envelope of the randomly changing absolute values of the amplitudes drops exponentially at $n \rightarrow \pm \infty$. As a
measure $\gamma$ of such localization one usually employs the limit

$$
\begin{equation*}
\gamma=\lim _{n \rightarrow \infty} \frac{1}{n}\left\langle\log \left(\psi_{n}^{2}+\psi_{n+1}^{2}\right)\right\rangle \tag{2.53}
\end{equation*}
$$

where the average $\langle\ldots\rangle$ is taken over the ensemble of possible realizations of the random potential, that is over $\left\{k_{n}\right\}$.

If the deviation $\delta k_{n}$ of each of the coefficients $k_{n}$ from the average is statistically independent of the others and has the probability

$$
\begin{equation*}
w\left(k_{n}\right)=\frac{1}{\pi} \frac{\delta}{\delta k_{n}^{2}+\delta_{n}^{2}} \tag{2.54}
\end{equation*}
$$

to have a given value $\delta k_{n}$, the problem has an exact solution which yields

$$
\begin{equation*}
\gamma=\operatorname{Arccosh}\left[\frac{1}{4} \sqrt{\left(2+\frac{E}{V}\right)^{2}+\frac{\delta^{2}}{V^{2}}}+\frac{1}{4} \sqrt{\left(2-\frac{E}{V}\right)^{2}+\frac{\delta^{2}}{V^{2}}}\right] \tag{2.55}
\end{equation*}
$$

The localization length is $L=1 / \gamma \simeq 2 V / \delta$ for small $E / V$ and $\delta / V$.

## Quantum States at a Surface

The localization phenomenon is considered as a universal property of onedimensional disordered quantum systems translationally invariant "on average". The situation may change when we are dealing with a set of level manifolds interacting in a relay-like manner. In this case one more dimension is involved in the problem such as the internal degrees of freedom of particles moving in a spatially limited domain. In Fig. 2.19(b) we illustrate this by the example of a particle which can move in a domain limited by a binding potential of a complex form. It can also tunnel to a similar potential of the neighboring domain.

The Schrödinger equation for this system reads

$$
\begin{equation*}
i \dot{\psi}_{n, k}=E_{n, k} \psi_{n, k}+\sum_{m}\left(V_{n, k}^{n+1, m} \psi_{n+1, m}+V_{n, k}^{n-1, m} \psi_{n-1, m}\right) \tag{2.56}
\end{equation*}
$$

where $V_{n, k}^{n \pm 1, m}$ are the tunneling amplitudes. One sees that this equation has the same form as (2.16) describing the excitation of a polyatomic molecule by an intense laser field. For random potential pits the spectra $E_{n, k}$ are random, and the matrix elements $V_{n, k}^{l, m}$ may also be random for the multidimensional internal motion. The latter is also the case for a particle on a random surface, for which the potential energy resembles a set of conducting fractals as shown in Fig. 2.20. For this case the problem of localization requires special consideration. Examples are known where, in two-dimensional disordered systems uniform "on average", the so-called weak localization phenomenon takes place: the particle does go to infinity, but it takes an exponentially long time.


Fig. 2.20. The random potential of a surface electron formed as a result of the deposition of conducting material onto an isolating surface (a). Tunneling interaction couples the one-electron states of different conducting ramified islands at the points $A, B, C$, and $D$. (b) The corresponding levels scheme.

## A classical Mechanics Analog

It is worth mentioning that one encounters equations of the type (2.52) or (2.56) not only for quantum motion of electrons, but also for classical vibrations of atoms in the lattice, namely phonons. This is the case when atoms of different masses are randomly distributed among the nodes.

After Fourier transformation of the Newton equation

$$
\begin{equation*}
m_{n} \ddot{x}_{n}=\frac{\partial U_{n}^{n-1}(d)}{\partial d}\left(x_{n}-x_{n-1}\right)-\frac{\partial U_{n+1}^{n}(d)}{\partial d}\left(x_{n}-x_{n-1}\right) \tag{2.57}
\end{equation*}
$$

for classical coordinates $x_{n}$ of the atoms of masses $m_{n}$, where $U_{n}^{n+1}(d)$ is the interaction potential between the $n$-th and $n+1$-th atoms at distance $d$, we arrive at

$$
\begin{equation*}
-\omega^{2} m_{n} x_{n}=\frac{\partial U_{n}^{n-1}(d)}{\partial d}\left(x_{n}-x_{n-1}\right)-\frac{\partial U^{n} n+1(d)}{\partial d}\left(x_{n}-x_{n-1}\right) . \tag{2.58}
\end{equation*}
$$

Equation (2.58) has the same structure as (2.52), since the randomness of masses plays the same role as the distribution of detunings $k_{n}$. For a given


Fig. 2.21. Oscillations in a chain of atoms with disordered masses. The envelope of the vibrational amplitudes of individual atoms drops exponentially.
vibrational mode the individual oscillation amplitudes of atoms change irregularly from one knot to another, but the envelope of this change drops exponentially as sketched in Fig. 2.21. This implies that phonons cannot propagate far in such a medium.

### 2.5.3 Nonlinear Optics

From the viewpoint of quantum mechanics many of the well-known phenomena described by the theory of nonlinear oscillations and nonlinear optics are processes in multilevel quantum systems. This fact usually does not play an important role in the limit of classical mechanics, that is when a process involves a large number of quanta. Nevertheless, it may be important for consideration of the threshold and subthreshold regimes of generation, chaotic regimes, or for the processes developing in small confined volumes and on long time-scales.

Let us illustrate this by the example of a nonlinear sample of volume $\mathcal{V}$ placed into a multimode resonator shown in Fig. 2.22. In the empty resonator only one mode, say $k_{0}$, is excited and contains $N$ photons, whereas all other modes $\left\{k_{n}\right\}$ are in the vacuum state. The electric field amplitude in the cavity $\mathcal{E}(\boldsymbol{r})=\sqrt{N} \mathcal{E}_{0}(\boldsymbol{r})$ is therefore given by the number of photons and the mode function $\mathcal{E}_{0}(\boldsymbol{r})$ of the excited mode. When we introduce the nonlinear medium to the resonator, the radiation of the mode $k_{0}$ induces a polarization $\mathcal{P}\left[\sqrt{N} \mathcal{E}_{0}(\boldsymbol{r})\right]$, which couples the excited mode to all other modes. The interaction energy

$$
\begin{equation*}
V_{n, 0}=-\int_{\mathcal{V}} d^{3} r \mathcal{P}\left[\sqrt{N} \mathcal{E}_{0}(\boldsymbol{r})\right] \mathcal{E}_{n}(\boldsymbol{r}) \tag{2.59}
\end{equation*}
$$

may change considerably and irregularly with the index $n$, if the sample is of an irregular shape, or if the nonlinear medium is close to saturation. In the


Fig. 2.22. Nonlinear interaction of modes in an optical resonator.
last case the irregularity arises from the spatial structure of the mode function $\mathcal{E}_{0}(\boldsymbol{r})$. We denote by $\psi_{0}$ the amplitude of the state $|N, 0,0, \ldots, 0\rangle$ with $N$ photons in the mode $k_{0}$ and no photons in other modes. By $\psi_{n}$ we denote the amplitude of the state $|N-1,0, \ldots, 1, \ldots, 0\rangle$ where one photon has been transported from the mode $k_{0}$ to the mode $k_{n}$. We neglect the probability of transportation of more than one photon and write the Schrödinger equation

$$
\begin{align*}
& i \hbar \dot{\psi}_{0}=\sum_{n} V_{0, n} \psi_{n} \\
& i \hbar \dot{\psi}_{n}=\Delta_{n} \psi_{n}+V_{m, 0} \psi_{0} \tag{2.60}
\end{align*}
$$

by analogy to (2.19). Here $\Delta_{n}$ is the difference in photon energies of the modes $k_{0}$ and $k_{n}$.

We mention here one more analogy. A nonlinear medium plays the role of an anharmonicity for the field oscillators in the resonator, similar to the role of vibrational anharmonicity for nuclear motion in polyatomic molecules, whence (2.60) describes the first quantum transition in such a system.

### 2.5.4 Cooperative Effect

Cooperative emission of an ensemble of non-interacting two-level atoms located in a region smaller than the wavelength of the resonant radiation is known as the Dicke effect. What happens if the atoms interact with each
other? Randomness in the mutual distances between the atoms and corresponding pair interactions suggests the random matrix approach.

The following analogy between the ensemble of $N$ two-level atoms and an ensemble of $N$ particles of spin $1 / 2$ in a magnetic field suggests a convenient means of description. One can consider each of the two-level atoms as a particle of spin $1 / 2$ in the magnetic field $\boldsymbol{B}$. Transition among the upper and the lower states of an atom corresponds to a spin flip from the state parallel to $\boldsymbol{B}$ to the antiparallel state, whereas the dipole-dipole interaction of atoms correspond to the spin-spin interaction.

The total spin $S$ of the ensemble may take any value from $N / 2$ to 0 for $N$ even, or to $1 / 2$ for $N$ odd. Each value of the total spin apart from $N / 2$ can be constructed in many ways by combining the directions of the spins of the individual particles. We mark each such combination by an index $\nu$. The statistical weight $\mathcal{G}_{S}$ of a state with spin $S$ is therefore given by the binomial coefficients $\mathcal{G}_{S}=\binom{S / 2}{N-S / 2}$ for $N$ even and $S$ integer, or $\mathcal{G}_{S}=\binom{S / 2-1 / 2}{N-S / 2+1 / 2}$ for $N$ odd and $S$ integer plus $1 / 2$. There are $2 S-1$ magnetic sublevels $M$ for any state of a given $S$. The energy of each sublevel depends linearly on $M$, and the spectrum of the ensemble is represented in terms of the quantum numbers $S$ and $M$, as shown in Fig. 2.23. Each energy level corresponds to a certain value of $M$, and even for a fixed $S$ is $\mathcal{G}_{S}$-fold degenerate with respect to $\nu$. Optical transitions in the system occur among the states with the same "spin" $S$ and the same number $\nu$.


Fig. 2.23. Energy spectrum of an ensemble of two-level systems in the representation of the total "spin" $S$ and its projection $M$.

Interaction of atoms changes considerably the structure of the spectrum, since it implies the possibility of transitions among all $\mathcal{G}_{S^{-}}$-states $|\nu\rangle$ with the same $S$ and $M$ values. The random arrangement of atoms results in a random coupling and yields a splitting of the degenerate levels. The new energy eigenstates $|k\rangle$ are linear combinations $|k\rangle=\sum s_{k, M}^{\nu, M}|\nu, M\rangle$. The radiative transitions induced by the electromagnetic field of amplitude $\mathcal{E}$ now occur in a band-like spectral structure, and therefore conform to the same model (see (2.16)) as anharmonic polyatomic molecules in a resonant laser field. The corresponding Schrödinger equation reads

$$
\begin{align*}
i \hbar \dot{\psi}_{M, k}= & \Delta_{M, k} \psi_{M, k}+\mathcal{E} d \sqrt{(M-S)(M+S+1)} \sum_{k^{\prime}} c_{M, k}^{k^{\prime}, M+1} \psi_{M+1, k^{\prime}} \\
& +\mathcal{E} d \sqrt{(M+S)(S+1-M)} \sum_{k^{\prime}} c_{M, k}^{k^{\prime}, M-1} \psi_{M-1, k^{\prime}} \tag{2.61}
\end{align*}
$$

where $\Delta_{M, k}$ are the interaction-induced random detunings of states from their original positions, $c_{k, M}^{k^{\prime}, M^{\prime}}=\sum_{\nu} s_{k, M}^{\nu, M}\left(s^{*}\right)_{k, M^{\prime}}^{\nu, M^{\prime}}$, and $d$ is the dipole moment of each atom.

### 2.5.5 Many-Body Effects in Cold Rydberg Gas

Consider an ensemble of cold atoms excited in a high Rydberg state $|n p\rangle$. Due to the dipole-dipole interaction, atoms may change their Rydberg state as a result of energy exchange: a couple of neighboring atoms initially in the state $|n p, n p\rangle$ performs a transition to the state $\left|n s, n^{\prime} s\right\rangle$. This energy exchange process may be viewed in the same way as the vibrational exchange in molecules described in Section 2.2: one atom emits a photon and performs a transition from the $|n p\rangle$ Rydberg state to the $|n s\rangle$ Rydberg state. The photon, being absorbed by the other atom, provokes a transition from the $|n p\rangle$ to the $\left|n^{\prime} s\right\rangle$ state (with $n^{\prime}=n+1$ ). Atomic transitions $n p \rightarrow n s$ and $n p \rightarrow n^{\prime} s$ are not always strictly in resonance, but they may be set to resonance by the external static electric field, which controls by the Stark effect the detuning $\Delta_{0}=\left(E_{n p}-E_{n s}\right)-\left(E_{n p}-E_{n^{\prime} s}\right)$ of the transition.

The process starts with pairs of relatively closed atoms and involves other atoms via the excitation exchange mechanism, as shown in Fig. 2.24(a). Each of the atoms of the pair can return to the state $|n p\rangle$ whereas two surrounding atoms perform the transitions to the $|n s\rangle$ and $\left|n^{\prime} s\right\rangle$ states. This resembles an autocatalytic process, where the products of a chemical reaction produced in a density fluctuation diffuse out of the reaction zone.

Note that the dipole moments of Rydberg atoms are large, and therefore the process of energy and excitation exchange is rapid for a reasonably high density $n_{a t}$ of the atomic vapor. It is much faster than the thermal motion of the cold atoms, and therefore the energy and excitation exchange in the ensemble can be considered as a Hamiltonian process without relaxation, in contrast to collisional relaxation in ensembles of warm atoms and molecules.

b)


Fig. 2.24. (a) A couple of $s$-states ( $s s^{\prime}$-couple) is created at a pair of close atoms (A and B). Excitation exchange processes result in the diffusion of these states out the pair. When the $s$-states have left the pair can produce the next $s s^{\prime}$-couple. (b) Level scheme of the elementary act. The initial state $|n p, n p\rangle$ corresponds to both atoms A and B of the pair in the state $|n p\rangle$. It is coupled by the energy exchange matrix element $V_{01}$ to the state where one of the atoms is in the Rydberg state $|n s\rangle$, while the other is in the state $|(n+1) s\rangle$. Diffusion of the $|n s\rangle$ and $|(n+1) s\rangle$ states over the ensemble of other atoms via the excitation exchange mechanism is represented by transitions to the band.

Viewing this process by analogy to the quasiparticle concept one can say that the couple of states $|n s\rangle$ and $\left|n^{\prime} s\right\rangle$ moves coherently to other sites as a result of the excitation exchange process. We refer to the $|n s\rangle$ and $\left|n^{\prime} s\right\rangle$ states as the $s$ and $s^{\prime}$ states respectively and to the $|n p\rangle$ state as the $p$ state. The excitation exchange results in the motion of the $s$-state over the immobile atoms, thus leading to diffusion of these states out of the "active center", that is, the pair of close atoms considered. Therefore, the pair of atoms can produce the next $s-s^{\prime}$ couple of states.

The level scheme for the elementary act of creation of a single $s-s^{\prime}$ couple is shown in Fig. 2.24(b). The initial state $|0\rangle$ corresponds to all atoms in the $|n p\rangle$ state. The first state $|1\rangle$ corresponds to atoms of a close pair (say the atom A and the atom B) in the states $|n s\rangle$ and $\left|n^{\prime} s\right\rangle$. The states $|0\rangle$ and $|1\rangle$ are coupled by the matrix element of the dipole-dipole interaction $\hat{V}_{1,2}=\mu_{n p}^{n s} \mu_{n p}^{n^{\prime} s} / R_{A B}^{3}$ responsible for the energy exchange process. The band corresponds to locations of $|n s\rangle$ and $\left|n^{\prime} s\right\rangle$ states on other atoms (say $\mathrm{C}, \mathrm{D}, \mathrm{F}$, etc.). The state $|1\rangle$ is coupled by the interaction to a combina-
tion $|b\rangle=\sum V_{1 n} /\left\langle V^{2}\right\rangle^{1 / 2}|n\rangle$, which is a properly normalized superposition of the band states. This combination represents the $s-s^{\prime}$ couple delocalized over a group of atoms near the initial pair A-B. It is not an eigenstate of the Hamiltonian, and hence is also coupled to another combination of the band states, etc. This hierarchy of couplings of the states corresponds to a dynamical process, which resembles diffusion with a diffusion coefficient $\simeq 4 \pi \mu^{2} n_{a t} \ln \left(R_{g} / a\right)$ where $R_{g}$ is the size of the Rydberg ensemble of the density $n_{a t} \sim 1 / a^{3}$, and $\mu$ is the dipole moment.

The Schrödinger equation reads

$$
\begin{align*}
i \dot{\psi}_{0} & =\Delta_{0} \psi_{0}+V_{01} \psi_{1}  \tag{2.62}\\
i \dot{\psi}_{1} & =V_{10} \psi_{0}+\sum_{n} V_{1 n} \psi_{n} \\
i \dot{\psi}_{n} & =\Delta_{n} \psi_{n}+V_{n 1} \psi_{1}
\end{align*}
$$

where $\psi_{0}, \psi_{1}$ and $\psi_{n}$ are the amplitudes of the states $|0\rangle,|1\rangle$ and $|n\rangle$ respectively, and $\Delta_{n}$ is the detuning of the state $|n\rangle$ in the band from the state $|1\rangle$. Analytical calculation of the energies $\Delta_{n}$ of the eigenstates and the coupling $V_{1 n}$ is a challenging task. Indeed, the mean interaction $\int_{0}^{\infty} V(R) n_{a t} d^{3} R$ in a gas with $1 / R^{3}$ binary interaction diverges both at small and long distances. This can be interpreted that far atoms and close atoms are equally important for the excitation-exchange process. In such a situation, the spatial structure of stationary states of the quasiparticles is not immediately evident - the stationary states may be localized on one or several atoms, or, on the contrary, each energy eigenstate may extend over the entire volume of the system, like in a solid. It turns out that both well-localized states and strongly delocalized states with multifractal spatial structure are possible, which yields a non-trivial profile of the level detuning statistics $g\left(\Delta_{n}\right)$. On the other hand, the coupling statistics $g\left(V_{1 n}\right)$ s imilar to that shown in Fig. 2.5 turn out to be statistically independent from detunings and to follow the inverse square law $g(V) \sim V^{-2}$.

## 3 Two-Level and Level-Band Systems

We start to consider the dynamics of populations of complex multilevel systems with two well-known examples: a two-level system and a system consisting of a level and a band which is either continuous or composed of discrete levels. The main purpose of such a consideration is to introduce the mathematical methods which are usually employed for the description, and to relate these methods to the main physical processes taking place in the systems, such as Rabi oscillations in two-level systems, repulsion of the levels by interaction, and the exponential decay of the level into the band, as well as to figure out what are the main physical parameters governing the population dynamics. The influence of the continuous band edges and the number of resonant states for the discrete band are discussed along with the process of population injection. We take into account the effects of interference resulting from the structure induced in the continuum by the presence of another interacting level and show how the interference yields special spectral structures known as Fano profiles.

For the discrete bands, we pay great deal of attention to the asymptotic behavior in the time domain longer than the Heisenberg return time. In this regime, the population returns from the band to the level. In a number of cases, the return can be complete and is then termed a revival. An interesting case of fractional revivals, when the population returns in the form of several short wavepackets, each of which contains a fraction of the complete population, is discussed in detail. The existence of fractional revivals suggests the revision of the Erenfest classical limit of quantum mechanics and its extension to the case where, in this limit, one can observe a quantum mixture of two or more classical wavepackets, the situation known as Schrödinger's cats.

We also concentrate on the case where an infinitely dense band has such an irregular structure that it cannot be considered as a continuum. Though this problem has an exact analytical solution in the form of a one-fold integral from an algebraic expression containing the sum of contributions of the individual band levels, the dynamics of such a system is not trivial. It consists of multiple returns of the population to the level that occur with different return times, differing by orders of magnitude. Interference of these returns results in a non-exponential decay of the level population averaged over a time interval. Moreover under certain conditions, the decay of the level can
be incomplete, which is the phenomenon related to dynamical localization. We find the parameters of the quantum system governing the population dynamics in this case.

### 3.1 Two-Level System

Let us concentrate first at a two-level system with the Hamiltonian

$$
\widehat{H}=\left(\begin{array}{cc}
E_{1} & 0  \tag{3.1}\\
0 & E_{0}
\end{array}\right)
$$

and write down the Schrödinger equation for the amplitudes $\Psi_{0}$ of the ground state, corresponding to the energy $E_{0}$ and the state vector $|0\rangle$. The amplitude, the energy, and the state vector of the exited state are $\Psi_{1}, E_{1}$, and $|1\rangle$ respectively. The Schrödinger equation for the Hamiltonian (3.1) reads

$$
\begin{align*}
i \hbar \dot{\Psi}_{1} & =E_{1} \Psi_{1} \\
i \hbar \dot{\Psi}_{2} & =E_{2} \Psi_{2} \tag{3.2}
\end{align*}
$$

What happens when we switch on an external periodic electromagnetic field of amplitude $\mathcal{E}$ and frequency $\omega$ close to the transition frequency $\omega \sim\left(E_{1}-\right.$ $\left.E_{0}\right) / \hbar$ if the interaction Hamiltonian is given by the product $\mathcal{E}_{\omega} \widehat{d}$ of the field strength $\mathcal{E}_{\omega}=\mathcal{E} \cos \omega t$ and the transition dipole moment operator $\widehat{d}$ ? We assume that the dipole moment operator has real non-diagonal matrix elements $d=\langle 0| \widehat{d}|1\rangle$, and the diagonal matrix elements vanish, as is usually the case in optics. The Schrödinger equation now takes the form

$$
\begin{align*}
& i \hbar \dot{\Psi}_{1}=E_{1} \Psi_{1}+d \mathcal{E} \cos (\omega t) \Psi_{0} \\
& i \hbar \dot{\Psi}_{0}=E_{0} \Psi_{0}+d \mathcal{E} \cos (\omega t) \Psi_{1} \tag{3.3}
\end{align*}
$$

We now apply to (3.3) the so-called rotating wave approximation, which is convenient when the detuning $\Delta=\left(E_{1}-E_{0}\right) / \hbar-\omega$ of the field frequency from the transition frequency $\left(E_{1}-E_{0}\right) / \hbar$ is relatively small $\Delta \ll \omega$, and which implies that the quantum phases of the states are taken relative to the phase $\omega t$ of the field. To this end we substitute the amplitudes of the ground and the exited states in the form

$$
\begin{align*}
\Psi_{1} & =e^{-i \omega t} \psi_{1} \\
\Psi_{0} & =\psi_{0} \tag{3.4}
\end{align*}
$$

to (3.3) and make use of the formula $\cos (\omega t)=\frac{1}{2}\left(e^{i \omega t}+e^{-i \omega t}\right)$, which yields

$$
\begin{align*}
i \hbar \dot{\psi}_{1} e^{-i \omega t}+\hbar \omega \psi_{1} e^{-i \omega t} & =E_{1} \psi_{1} e^{-i \omega t}+\frac{d \mathcal{E}}{2}\left(e^{i \omega t}+e^{-i \omega t}\right) \psi_{0} \\
i \hbar \dot{\psi}_{0} & =E_{0} \psi_{0}+\frac{d \mathcal{E}}{2}\left(e^{i \omega t}+e^{-i \omega t}\right) \psi_{1} e^{-i \omega t} \tag{3.5}
\end{align*}
$$

When we set $E_{0}=0$ and replace the combination $d \mathcal{E} / 2$ by a shorter notation $V$, after multiplication of the first equation by $e^{i \omega t}$ the system (3.5) takes the form

$$
\begin{align*}
& i \hbar \dot{\psi}_{1}=\left(E_{1}-\hbar \omega\right) \psi_{1}+V\left(1+e^{2 i \omega t}\right) \psi_{0} \\
& i \hbar \dot{\psi}_{0}=V\left(1+e^{-2 i \omega t}\right) \psi_{1} \tag{3.6}
\end{align*}
$$

To shorten notation further we also set $\hbar=1$ hereafter.
We now note that typical time derivatives of the amplitudes $\psi_{1}$ and $\psi_{0}$ in (3.6) are either of the order of the size of the interaction $V$ or of the order of the detuning $\Delta=\left(E_{1}-\omega\right)$, when it exceeds $V$, and therefore these derivatives are much smaller compared to the frequency of the external field $\omega$, provided the condition of resonance $\Delta \ll \omega$ is fulfilled. This enables us to ignore the rapidly oscillating terms $V e^{2 i \omega t}$ in (3.6) that give negligibly small average contributions, and we arrive at

$$
\begin{align*}
i \dot{\psi}_{1} & =\Delta \psi_{1}+V \psi_{0} \\
i \dot{\psi}_{0} & =V \psi_{1} \tag{3.7}
\end{align*}
$$

Let us now consider the solution of (3.7) corresponding to the initial conditions

$$
\begin{align*}
& \psi_{0}(t=0)=1 \\
& \psi_{1}(t=0)=0 . \tag{3.8}
\end{align*}
$$

One of our main tools is a method known either as the Fourier-Laplace transformation or as the Fourier transformation of generalized functions or distributions. The essence of the method is the following. The regular technique of Fourier transformations requires that the subjecting function vanishes at $t \rightarrow \pm \infty$, otherwise the Fourier transform of a derivative $\dot{y}(\omega)$, denoted by $F[\dot{y}(t)]$, differs from the Fourier transform of the primitive function multiplied by the frequency $\dot{y}(\omega) \neq i \omega y(t)$. Therefore all the functions $\psi$ are assumed to be zero at $t<0$ and coincide with the solutions of (3.7) with the initial conditions (3.8) for $t>0$. This implies that the function $\psi_{0}(t)$ is discontinuous at $t=0$ and it experiences a jump from the value $\psi_{0}=0$ to the value $\psi_{0}=1$ at this point. Apparently, the derivative $\dot{\psi}_{0}$ equals infinity, or in other words it amounts to the Dirac $\delta$-function. Equation (3.7) is no longer valid at $t=0$, and the correct equation reads

$$
\begin{align*}
i \dot{\psi}_{1} & =\Delta \psi_{1}+V \psi_{0} \\
i \dot{\psi}_{0} & =V \psi_{1}+i \delta(t) \tag{3.9}
\end{align*}
$$

We note here, that the Dirac $\delta$-function in the second equation (3.9) can be interpreted in physical terms. Since we assume that at $t<0$ there was no particle in the system of two levels, and at $t=0$ a particle has been supplied
to the ground state $|0\rangle$, one can say that the flux of the probability amplitude of the particle supply is infinitely large during an infinitely short time and is given by the $\delta$-function. Such a situation suggests an evident generalization: one can supply a particle with the $\Pi=\Pi(t)$, which can be an arbitrary function of time satisfying just a normalization condition. Such a situation is shown schematically in Fig. 3.1


Fig. 3.1. Two-level system in an external resonant field. A quantum particle is supplied at the lower state with the probability amplitude flux $\Pi(t)$. A transition dipole moment $d$ results in the coupling $V=\mathcal{E} d$ between the states. (a) Representation in the basis of energy eigenstates. (b) Representation in the rotating wave approximation.

Let us now perform the Fourier transformation of (3.9) and arrive at

$$
\begin{align*}
\varepsilon \psi_{1}(\varepsilon)-\Delta \psi_{1}(\varepsilon)-V \psi_{0}(\varepsilon) & =0 \\
\varepsilon \psi_{0}(\varepsilon)-V \psi_{1}(\varepsilon) & =i \tag{3.10}
\end{align*}
$$

where

$$
\begin{equation*}
\psi_{0 ; 1}(\varepsilon)=\int_{-\infty}^{+\infty} e^{i \varepsilon t} \psi_{0 ; 1}(t) d t \tag{3.11}
\end{equation*}
$$

and we have employed the fact that $\int_{-\infty}^{+\infty} e^{i \varepsilon t} \delta(t) d t=1$.
The set of two linear algebraic equations (3.10) has the solution

$$
\begin{align*}
& \psi_{1}(\varepsilon)=\frac{\left|\begin{array}{cc}
0, & -V \\
i, & \varepsilon
\end{array}\right|}{\left|\begin{array}{cc}
\varepsilon-\Delta, & -V \\
-V, & \varepsilon
\end{array}\right|}=\frac{V i}{(\varepsilon-\Delta) \varepsilon-V^{2}} \\
& \psi_{0}(\varepsilon)=\frac{\left|\begin{array}{cc}
\varepsilon-\Delta, & 0 \\
-V, & i
\end{array}\right|}{\left|\begin{array}{cc}
\varepsilon-\Delta, & -V \\
-V, & \varepsilon
\end{array}\right|}=\frac{(\varepsilon-\Delta) i}{(\varepsilon-\Delta) \varepsilon-V^{2}}=\frac{i}{\varepsilon-\frac{V^{2}}{(\varepsilon-\Delta)}} \tag{3.12}
\end{align*}
$$

where vertical brackets around a matrix denote the determinants. We find the time-dependent probability amplitudes $\psi_{0}(t)$ and $\psi_{1}(t)$ by performing the inverse Fourier transformation

$$
\begin{equation*}
\psi_{0 ; 1}(t)=\frac{1}{2 \pi} \int_{-\infty+i \nu}^{+\infty+i \nu} e^{-i \varepsilon t} \psi_{0 ; 1}(\varepsilon) d(\varepsilon) \tag{3.13}
\end{equation*}
$$

that is by integrating the Fourier transforms $\psi_{0}(\varepsilon)$ and $\psi_{1}(\varepsilon)$ in the complex plane along the contour $C$ given by a straight line which goes from $-\infty$ to $+\infty$ at a distance $\nu$ above the real axis, as shown in Fig. 3.2


Fig. 3.2. The integration contour of the inverse Fourier transform goes above the real axis. It can be closed in the lower part of the complex plane of $\varepsilon$, where the integrand vanishes.

The denominators of the right-hand sides of (3.12) tend to zero at the points

$$
\begin{equation*}
\varepsilon_{0 ; 1}=\frac{\Delta}{2} \pm \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}} \tag{3.14}
\end{equation*}
$$

we also note that the integrals (3.13) taken along the contour $C_{2}$ (a circle of large radius $\varepsilon=R \exp (-i \theta) ; R \longrightarrow \infty ; 0 \leq \theta \leq \pi)$ vanish in the limit
$R \longrightarrow \infty$ since the exponent of the integral acquires a large negative part. We can therefore replace the open integration contour $C_{1}=(-\infty-i \nu, \infty-i \nu)$ by the closed integration contour $C=C_{1}+C_{2}$ and make use of the Cauchy theorem by replacing the integral by the sum of the contributions of residuals corresponding to the singularities of the integrand at the points $\varepsilon_{0 ; 1}(3.14)$ in the lower part of the complex plane $\varepsilon$ including the real axis. We arrive at

$$
\begin{align*}
& \psi_{1}(t)=\frac{V}{\left(\varepsilon_{0}-\varepsilon_{1}\right)} e^{-i \varepsilon_{0} t}-\frac{V}{\left(\varepsilon_{0}-\varepsilon_{1}\right)} e^{-i \varepsilon_{1} t} \\
& \psi_{0}(t)=\frac{\left(\varepsilon_{0}-\Delta\right) V}{\left(\varepsilon_{0}-\varepsilon_{1}\right)} e^{-i \varepsilon_{0} t}-\frac{\left(\varepsilon_{1}-\Delta\right) V}{\left(\varepsilon_{0}-\varepsilon_{1}\right)} e^{-i \varepsilon_{1} t} \tag{3.15}
\end{align*}
$$

which yields

$$
\begin{align*}
\psi_{1}(t)= & \frac{V}{2 \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}} e^{-i\left(\frac{\Delta}{2}+\sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}\right) t}-\frac{V}{2 \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}} e^{-i\left(\frac{\Delta}{2}-\sqrt{\left.\left(\frac{\Delta}{2}\right)^{2}+V^{2}\right) t}\right.} \\
\psi_{0}(t)= & \frac{\left(\frac{\Delta}{2}+\sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}\right)}{2 \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}} e^{-i\left(\frac{\Delta}{2}+\sqrt{\left.\left(\frac{\Delta}{2}\right)^{2}+V^{2}\right)}\right) t}  \tag{3.16}\\
& -\frac{\left(\frac{\Delta}{2}-\sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}\right)}{2 \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}} e^{-i\left(\frac{\Delta}{2}-\sqrt{\left.\left(\frac{\Delta}{2}\right)^{2}+V^{2}\right) t}\right.}
\end{align*}
$$

that is

$$
\begin{align*}
& \psi_{1}(t)=-\frac{V i}{\sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}} e^{-i \frac{\Delta}{2} t} \sin t \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}} \\
& \psi_{0}(t)=e^{-i \frac{\Delta}{2} t}\left[\cos t \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}-\frac{\frac{\Delta}{2} i \sin t \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}}{\sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}}\right] . \tag{3.17}
\end{align*}
$$

The combination $\sqrt{(\Delta / 2)^{2}+V^{2}}$ is known as the Rabi frequency, and the corresponding periodic oscillations of the probabilities are referred to as the Rabi oscillations.

Let us now calculate the populations $\rho_{1}=\left|\psi_{1}\right|^{2}$ of the upper and $\rho_{0}=$ $\left|\psi_{0}\right|^{2}$ of the lower levels. Equation (3.17) immediately results in

$$
\begin{align*}
& \rho_{1}(t)=\frac{V^{2}}{\left(\frac{\Delta}{2}\right)^{2}+V^{2}} \sin ^{2} t \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}} \\
& \rho_{0}(t)=\frac{\Delta^{2}}{\Delta^{2}+4 V^{2}} \sin ^{2} t \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}+\cos ^{2} t \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}} . \tag{3.18}
\end{align*}
$$

Let us now ask for the time average populations $\overline{\rho_{0 ; 1}}$ of the states. Allowing for $\overline{\cos ^{2} \Omega t}=\overline{\sin ^{2} \Omega t}=1 / 2$ we arrive at

$$
\begin{align*}
& \overline{\rho_{1}}=\frac{2 V^{2}}{\Delta^{2}+4 V^{2}} \\
& \overline{\rho_{0}}=\frac{1}{2}\left[\frac{\Delta^{2}}{\Delta^{2}+4 V^{2}}+1\right] \tag{3.19}
\end{align*}
$$

One can see that for $V \ll \Delta / 2$ the mean population $\rho_{1}$ of the upper level is small, and the population $\rho_{0}$ is of the order of 1 . The populations of both levels become of the same order of magnitude only when $V \sim \Delta / 2$. At the extreme $V \gg \Delta / 2$ we have $\rho_{1} \simeq \rho_{0} \simeq 1 / 2$. We can therefore say that the states $|0\rangle$ and $|1\rangle$ are in resonance and the condition of the resonance $V \sim \Delta / 2$ is that the typical interaction amplitude becomes larger or of the order of the typical detuning. Later on we will often employ this qualitative criterion analyzing the behavior of multilevel quantum systems.

We also note an important case when the detuning is so small that it can be set zero. In this case

$$
\begin{align*}
& \rho_{1}(t)=\sin ^{2} V t \\
& \rho_{0}(t)=\cos ^{2} V t \tag{3.20}
\end{align*}
$$

and the population oscillates between the lower and the upper states with the frequency $V$. In this case the average populations of the states are apparently equal and amount to $1 / 2$.

Sometimes it is expedient to represent the time-dependent amplitudes $\psi_{0,1}(t)(3.17)$ in the form of a superposition of exact energy eigenstates $|u\rangle$, $|l\rangle$ of the system corresponding to the upper and lower energy eigenvalues $\varepsilon_{1 ; 0}=\frac{\Delta}{2} \pm \sqrt{\left(\frac{\Delta}{2}\right)^{2}+V^{2}}$ of (3.14) respectively. The amplitudes $\varphi_{l, u}$ of these states have only the phase time dependence $e^{-i \varepsilon_{0 ; 1} t}$ whereas the absolute values of the amplitudes and the phase factors at $t=0$ one finds with the help of the representation

$$
\begin{align*}
& |l\rangle=\frac{|0\rangle\left(\frac{\Delta}{2 V}+\sqrt{1+\frac{\Delta^{2}}{4 V^{2}}}\right)}{\sqrt{1+\left(\frac{\Delta}{2 V}+\sqrt{1+\frac{\Delta^{2}}{4 V^{2}}}\right)^{2}}+\frac{|1\rangle}{\sqrt{1+\left(\frac{\Delta}{2 V}+\sqrt{\left.1+\frac{\Delta^{2}}{4 V^{2}}\right)^{2}}\right.}}} \begin{array}{|l|}
|u\rangle=\frac{|0\rangle \sqrt{1+\left(\frac{\Delta}{2 V}-\sqrt{1+\frac{\Delta^{2}}{4 V^{2}}}\right)^{2}}}{2 \sqrt{1+\frac{\Delta^{2}}{4 V^{2}}}}-\frac{|1\rangle}{\sqrt{1+\left(\frac{\Delta}{2 V}-\sqrt{1+\frac{\Delta^{2}}{4 V^{2}}}\right)^{2}}}
\end{array} .=\begin{array}{l}
\sqrt{\sqrt{2}}
\end{array} \tag{3.21}
\end{align*}
$$

of the eigenvectors in the basis of states $|0\rangle$ and $|1\rangle$. In particular, for the case $\Delta=0$ we have two eigenstates

$$
\begin{align*}
& |l\rangle=|+\rangle=\frac{1}{\sqrt{2}}|0\rangle+\frac{1}{\sqrt{2}}|1\rangle \\
& |u\rangle=|-\rangle=\frac{1}{\sqrt{2}}|0\rangle-\frac{1}{\sqrt{2}}|1\rangle \tag{3.22}
\end{align*}
$$

corresponding to the eigenenergies $\varepsilon_{0 ; 1}= \pm V$, with the relative phase difference $2 V t$ of the amplitudes $\varphi_{u, l}=e^{ \pm i V t} / \sqrt{2}$ having identical initial values $\langle l \mid 0\rangle=\langle u \mid 0\rangle=1 / \sqrt{2}$. Beating of the oscillations of these two amplitudes yields (3.20) for the populations.

### 3.2 Level-band System

We continue by tackling a more complex system of a level interacting with a band of levels. We first consider it in the general case and obtain an exact expression for the Fourier transforms of time-dependent amplitudes of the level population amplitudes by the analogy to the results of the previous section. It turns out that this problem possesses an exact analytical solution. However, physically meaningful results emerge from the general expressions only in several limiting cases. The most well-known limit corresponds to an infinite density of states and a constant interaction and yields an expression known as the Fermi golden rule. The other important case is known as the Fano problem. This also implies an infinite density of states but allows for a non-uniform, that is energy-dependent, amplitude of the transition from the initial level to the band which by itself is a result of the interaction of another discrete level with a continuum. Comparison of the solutions of these two problems shows an important role of the quantum interference that affects strongly the population dynamics. This role becomes even more important when we consider a band of finite state density, that is a zone composed by a large but finite number of discrete levels. It results in specific effects at long times-scale - recurrences, revivals, and fractional revivals. We also study the dynamics of the system in the case, when the coupling among the level and the states of the band can considerably, that is by orders of magnitudes, vary from state to state. We show that the interference changes completely the dynamics of the system in the long times-scale, although this phenomenon does not manifest itself in periodic oscillations typical of interference phenomena.

### 3.2.1 General Consideration

In the previous section by having employed the rotating wave representation and the resonance approximation we have shown that a system in a periodic external field behaves as a system subjected to a time-independent perturbation, provided the detuning of the external frequency from the transition frequency is relatively small. Considering here a level-band system we make
use of the same approximation and by the analogy arrive at the Shrödinger equation

$$
\begin{align*}
i \dot{\psi}_{n} & =\Delta_{n} \psi_{n}+V_{n}^{0} \psi_{0} \\
i \dot{\psi}_{0} & =\sum_{n=1}^{N} V_{0}^{n} \psi_{n} \tag{3.23}
\end{align*}
$$

Here $\psi_{0}$ is a probability amplitude of the system to be in the ground state, and $\psi_{n}$ denote the amplitudes of the energy eigenstates in the $N$-level band. The detunings $\Delta_{n}$ and the interaction amplitudes $V_{0}^{n}$ now depend on the numbers $n$ of the band levels that are supposed to be non-interacting among themselves. In Fig. 3.3 we illustrate the level scheme of the system.


Fig. 3.3. Level-band system. The levels $|n\rangle$ of the band detuned by $\Delta_{n}$ are not coupled among themselves directly, but only via couplings $V_{0}^{n}$ with the level $|0\rangle$. The system either was initially in the state $|0\rangle$ with probability amplitude $\psi_{0}(t=0)=1$, or has been injected into this state with probability amplitude flux $\Pi(t)$.

Let us assume that at $t=0$ only the state $\psi_{0}$ is populated. Then by the analogy to (3.9) we write

$$
\begin{align*}
i \dot{\psi}_{n} & =\Delta_{n} \psi_{n}+V_{n}^{0} \psi_{0} \\
i \dot{\psi}_{0} & =\sum_{n=1}^{N} V_{0}^{n} \psi_{n}+i \delta(t) \tag{3.24}
\end{align*}
$$

For the Fourier transforms $\psi_{n}(\varepsilon)$ by analogy to (3.10) we therefore have

$$
\begin{align*}
\varepsilon \psi_{n}-\Delta_{n} \psi_{n}-V_{n}^{0} \psi_{0} & =0 \\
\varepsilon \psi_{0}-\sum_{n=1}^{N} V_{0}^{n} \psi_{n} . & =i \tag{3.25}
\end{align*}
$$

We now express $\psi_{n}(\varepsilon)$ in terms of $\psi_{0}(\varepsilon)$ with the help of the first equation (3.25), substitute this expression to the second equation, factor out the common multiple $\psi_{0}(\varepsilon)$ from the sum, and finally arrive at

$$
\begin{align*}
\psi_{0}(\varepsilon) & =\frac{i}{\varepsilon-\sum_{n=1}^{N} \frac{V_{0}^{n} V_{n}^{0}}{\varepsilon-\Delta_{n}}} \\
\psi_{n}(\varepsilon) & =\frac{V_{0}^{n}}{\varepsilon-\Delta_{n}} \psi_{0}=\frac{V_{0}^{n}}{\varepsilon-\Delta_{n}} \frac{i}{\varepsilon-\sum_{n=1}^{N} \frac{V_{0}^{n} V_{n}^{0}}{\varepsilon-\Delta_{n}}} \tag{3.26}
\end{align*}
$$

which is a formal solution of the problem in the Fourier representation.

### 3.2.2 Continuous Band Model

Equations (3.26) give an exact solution of the problem for arbitrary amplitudes of interaction $V_{0}^{n}$ and arbitrary detunings $\Delta_{n}$. It allows for all possible interference phenomena that may occur in the system. However the practical usefulness of these equations is limited, since they contain a complicated function given by the sum in the denominators. Therefore evaluation of the integrals in the inverse Fourier transformation now becomes the main difficulty. Hence for the further analysis of the problem we have to find a reasonable approximation for the sum, which would allow us to find an analytical solution of the problem. One possibility is to replace these sums by an integral which can be evaluated. The following physical example of a level-band system will help us to gain an insight into the physical meaning of the approximation.

## An Example of the Level-Band Model

Let us consider a quantum particle moving in an external field formed by two potential wells $A$ and $B$ separated by a barrier $C$ shown in Fig. 3.4. The first well is deep and narrow and the second is wide and shallow. The distance between successive levels is of the order of the oscillation frequency, and therefore the spacing among the neighboring levels $|n\rangle$ in the shallow well $B$ is much smaller than that in the deep well $A$.

Here we restrict our consideration to the energy region around the first excited level in the first well, and allow for the tunneling from this state to the states of the potential well $B$ through the potential barrier $C$. We therefore assign index 0 to this level, which is the initial state of our level-band problem. The lower state in the pit $A$ serves only as a reservoir for the probability amplitude supply $\Pi(t)$, and will not be taken into account otherwise. Note, that energy eigenstates in each of the pits are not the exact eigenstates of the entire problem. However they can be taken as a good approximative basis set, provided the barrier is sufficiently high. These basis states are not orthogonal, and the Hamiltonian is not diagonal in such a representation. It contains offdiagonal matrix elements given by the overlap integrals, which are of the order of tunneling exponents $V_{0}^{n} \sim \exp \left\{-\operatorname{Im} \int_{\text {under barrier } C} \sqrt{E-U(x)} d x\right\}$. Here we take all $V_{0}^{n}$ to be identical, since the distances among the levels of interest are much smaller than the height of the barrier, and all tunneling exponents are almost equal.


Fig. 3.4. A quantum particle in a two-well potential, which consists of two pits $A$ and $B$ separated by a barrier $C$, is an example of a level-band system with an amplitude injection to the level. The interaction $V_{0}^{n}$ among the energy eigenstates of different pits is of the order of the tunneling exponent. The spacing $\delta$ of neighboring levels in the pit of length $L$ is of the order of the inverse period, that is the inverse time of flight from wall to wall. The population supply $\Pi(t)$ to the exited state in the left pit can be controlled by an external field.

Let us turn now to the spacing $\delta$ among the neighboring levels in the pit $B$. For a rectangular well the energy eigenfunctions satisfy the Bohr-Sommerfeld quantization condition

$$
\begin{align*}
& \text { right turning point } \\
& \qquad \int_{\text {left turning point }} \sqrt{2 m E_{n}} d x=\sqrt{2 m E_{n}} L=n \pi \tag{3.27}
\end{align*}
$$

From (3.27) one immediately finds the energy difference among the neighboring levels:

$$
\begin{equation*}
\sqrt{2 m E_{n}}-\sqrt{2 m E_{n-1}}=\sqrt{\frac{m}{2 E_{n}}}\left(E_{n}-E_{n-1}\right)=\frac{\pi}{L} . \tag{3.28}
\end{equation*}
$$

We note that $v=\sqrt{m / 2 E_{n}}$ is a typical velocity of the particle in the potential well $B$ corresponding to the energy region $E_{n}$, and therefore

$$
\begin{equation*}
\delta=\left(E_{n}-E_{n-1}\right)=\frac{\pi}{L} v \simeq \frac{\pi}{T} \tag{3.29}
\end{equation*}
$$

where $T$ is the time of flight from the left wall of the pit to the right wall. The level spacing therefore amounts to the reciprocal time of return flight in the pit.

We immediately understand the meaning of the continuum model limit $\delta \rightarrow 0$ : it means that the particle never returns once it has left the potential pit $A$, since the limit $T \rightarrow \infty$ means that the right turning point of the potential $B$ is infinitely far away from the barrier. We can phrase it differently by saying that one can replace a discrete spectrum by the continuum when considering the evolution during a time interval shorter than the inverse spacing among the neighboring levels, that is the state density $t \ll 1 / \delta \sim g$ or the Heisenberg return time. Note that it is exactly the same regime when the uncertainty of the energy due to the finite observation time exceeds the level spacing $\delta$.

## Level and Uniform Continuum

Let us now consider the consequences of the continuous-band approximation. In (3.26) we can now factor out the moduli squared of the interaction amplitudes $V^{2}=\left|V_{0}^{n}\right|^{2}$, substitute $\Delta_{n}=n \delta$, and replace summation over $n$ by the integration. We arrive at

$$
\begin{align*}
\psi_{0}(\varepsilon) & =\frac{i}{\varepsilon-V^{2} \int \frac{d n}{\varepsilon-n \delta}} \\
\psi_{n}(\varepsilon) & =\frac{V}{\varepsilon-n \delta} \frac{i}{\varepsilon-V^{2} \int \frac{d n}{\varepsilon-n \delta}} \tag{3.30}
\end{align*}
$$

We now remember Fig. 3.2 where the integration contour of the inverse Fourier transformation is shown. It goes above the real axis, which implies that in (3.30) $\varepsilon$ is a complex number with a positive imaginary part. We also note that the strongly detuned levels, corresponding to large positive and negative $n$, are so far that they should not make a physically important contribution to the process. Therefore we can extend our integration over $d n$ to all the real axis from $-\infty$ to $+\infty$. We calculate the integral and find

$$
\begin{equation*}
\int_{-\infty}^{+\infty} \frac{d n}{\varepsilon-n \delta}=-\frac{1}{\delta} \int_{-\infty}^{+\infty} \frac{d x}{x-\varepsilon / \delta}=-\frac{1}{\delta} \ln \frac{\infty+i \nu / \delta}{-\infty+i \nu / \delta}=-\frac{i \pi}{\delta} \tag{3.31}
\end{equation*}
$$

where we have employed the explicit expression $\log (z)=\log |z|+i \arg (z)$ for the logarithm of a complex number $z$. After substitution of (3.31) into (3.30) we perform the inverse Fourier transformation and obtain

$$
\begin{align*}
& \psi_{0}(t)=\frac{1}{2 \pi} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{i}{\varepsilon-i V^{2} \pi / \delta} e^{i \varepsilon t} d \varepsilon \\
& \psi_{n}(t)=\frac{1}{2 \pi} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{V}{\varepsilon-n \delta} \frac{i}{\varepsilon-i V^{2} \pi / \delta} e^{i \varepsilon t} d \varepsilon \tag{3.32}
\end{align*}
$$

We close the integration contour in the lower plane and apply the Cauchy formula, which yields

$$
\begin{align*}
\psi_{0}(t) & =e^{-t V^{2} \pi / \delta} \\
\psi_{n}(t) & =\frac{e^{i t n \delta}}{n \delta-i V^{2} \pi / \delta}+\frac{e^{-t V^{2} \pi / \delta}}{i V^{2} \pi / \delta-n \delta} \tag{3.33}
\end{align*}
$$



Fig. 3.5. Level-band problem. (a) The Lorentzian distribution of the population over the states of the continuum at the long time limit, when the density of states $g=1 / \delta$ increases and the coupling $V$ decreases, such that the product $W=V^{2} g$ remains constant. For the example of Fig. 3.4 such a distribution corresponds to a wavepacket (b) which propagates in the continuum to the right and has an exponential envelope.

One can see that the population $\rho_{0}(t)=\left|\psi_{0}\right|^{2}$ of the state $|0\rangle$ decays exponentially

$$
\begin{equation*}
\rho_{0}(t)=e^{-t 2 \pi V^{2} / \delta} \tag{3.34}
\end{equation*}
$$

with the rate $W=2 \pi V^{2} / \delta$. It implies that at the time $T \gg t \gg \delta / 2 \pi V^{2}$ all the population is on the levels $|n\rangle$ of the band. The second equation (3.33) allows one to find a distribution of the population over these levels. The second term of the equation vanishes at $t \rightarrow \infty$ and the first term yields

$$
\begin{equation*}
\rho_{n}(t=\infty)=\frac{V^{2}}{(n \delta)^{2}+\left(V^{2} \pi / \delta\right)^{2}} \tag{3.35}
\end{equation*}
$$

which means that the population is distributed over the energy levels of the band according to a Lorentzian profile, as shown in Fig. 3.5. Integration of the right-hand side of (3.35) over $n$ yields unity, as it should be in accordance with the normalization condition.

Let us now rewrite (3.35) in terms of the detuning $\Delta_{n}=n \delta$ of the state $|n\rangle$ from resonance with the state $|0\rangle$, omit the index $n$, since it enters the expression only via the detuning, replace $V^{2} \pi / \delta$ by $W$, divide both parts by $\delta$, and obtain the spectral density of the population:

$$
\begin{equation*}
f(\Delta)=\frac{\rho(\Delta, t=\infty)}{\delta}=\frac{1}{\pi} \frac{W}{\Delta^{2}+W^{2}} \tag{3.36}
\end{equation*}
$$

This quantity has the same Lorentzian form and the same physical origin as spectral line shapes of width $W$. Indeed, the Dirac $\delta$-like supply $\Pi(t)$ in (3.24) can be represented as a superposition of harmonics of all frequencies with equal amplitudes. The Schrödinger equation is linear, and hence each harmonic at a frequency $\Delta$ brings at $t=\infty$ an independent contribution to the overall population transfer at this frequency. The efficiency of this process depends on the absorption rate at this frequency given by the Lorentzian profile, and thereby the final result coincides with (3.36).

## Role of the Continuum Edges

Let us now consider the role of the continuum edges. In (3.30) we have to take into account finite integration limits that yield

$$
\begin{align*}
\psi_{0}(\varepsilon) & =\frac{i}{\varepsilon-V^{2} \int_{a}^{b} \frac{d n}{\varepsilon-n \delta}}=\frac{i}{\varepsilon+\frac{V^{2}}{\delta} \ln \left(\frac{-\varepsilon+b}{-\varepsilon+a}\right)} \\
\psi_{n}(\varepsilon) & =\frac{V}{\varepsilon-n \delta} \frac{i}{\varepsilon+\frac{V^{2}}{\delta} \ln \left(\frac{-\varepsilon+b}{-\varepsilon+a}\right)} \tag{3.37}
\end{align*}
$$

The time dependence of the probability amplitude $\psi_{0}(t)$ depends on the positions of poles in the complex plane, which are given by roots of the equation

$$
\begin{equation*}
\varepsilon+\frac{V^{2}}{\delta} \ln \left(\frac{-\varepsilon+b}{-\varepsilon+a}\right)=0 \tag{3.38}
\end{equation*}
$$

For $a<\operatorname{Re}[\varepsilon]<b$ the logarithm has an imaginary part and corresponding poles are not disposed on the real axis. Therefore they bring to the probability amplitudes contributions exponentially decreasing in time. Besides on the complex solutions, equation (3.38) has two real roots at the points $\varepsilon_{a}<a$ and $\varepsilon_{b}>b$ that bring two oscillating contributions. These non-vanishing contributions dominate at the asymptotic regime of long time even in the case their amplitudes are small. In Fig. 3.6 we compare two terms of (3.38), one of which is taken with the opposite sign. Crossing points of the curves correspond to the non-vanishing oscillating contributions

$$
\begin{align*}
\psi_{0}(t) & =-\frac{2 \pi\left(\varepsilon_{a}-b\right)\left(\varepsilon_{a}-a\right) e^{-i \varepsilon_{a} t}}{\left(\varepsilon_{a}-b\right)\left(\varepsilon_{a}-a\right)+\frac{V^{2}}{\delta}(b-a)}-\frac{2 \pi\left(\varepsilon_{b}-b\right)\left(\varepsilon_{b}-a\right) e^{-i \varepsilon_{b} t}}{\left(\varepsilon_{b}-b\right)\left(\varepsilon_{b}-a\right)+\frac{V^{2}}{\delta}(b-a)} \\
& \simeq-\frac{2 \pi \delta}{V^{2}}\left[\left(\varepsilon_{a}-a\right) e^{-i \varepsilon_{a} t}+\left(\varepsilon_{b}-b\right) e^{-i \varepsilon_{b} t}\right] \tag{3.39}
\end{align*}
$$

of the amplitudes $2 \pi \delta\left(\varepsilon_{a}-a\right) / V^{2}$ and $2 \pi \delta\left(\varepsilon_{b}-b\right) / V^{2}$ that are small provided the bandwidth $(b-a)$ considerably exceeds the Lorentzian line width


Fig. 3.6. Graphical solution of (3.38). For $\varepsilon<a$ and $\varepsilon>b$ two real roots correspond to oscillating non-vanishing contributions, and their amplitude equal the difference of the slopes at the crossing point.
$W=V^{2} / \delta$. Indeed, for the lower edge of the band $\varepsilon_{b} \sim b<0$ equation (3.38) yields $\delta\left(b-\varepsilon_{b}\right) / V^{2} \simeq e^{b \delta / V^{2}}(a-b) \delta / V^{2}$. By analogy $(\varepsilon-a) \delta / V^{2} \simeq$ $e^{-a \delta / V^{2}}(a-b) \delta / V^{2}$ implies $a>0$, and therefore

$$
\begin{equation*}
\psi_{0}(t) \simeq \frac{(a-b) \delta}{V^{2}}\left[e^{b \delta / V^{2}-i \varepsilon_{b} t}-e^{-a \delta / V^{2}-i \varepsilon_{a} t}\right] \tag{3.40}
\end{equation*}
$$

One sees that the time-independent amplitudes are exponentially small when the edges of the band are disposed at distances larger compared to the width $W$ of the Lorentzian (3.36).

## Number of Levels in Resonance

We note that the result (3.35) is in complete agreement with the reasoning based on the time-energy uncertainty principle: the inverse life-time $1 / \tau$ given by the decay rate $W$ determines the distribution width. It implies that the population is distributed among the $\mathcal{N} \simeq W / \delta \sim V^{2} / \delta^{2}$ states of the band closest to the resonance with the state $|0\rangle$. We also note that the ratio $V^{2} / \delta^{2}$ should be a large number, otherwise it is impossible to satisfy the condition that the decay time $\tau=\delta / V^{2}$ is shorter than the Heisenberg return time $1 / \delta$ in the zone, which has been employed for replacing of the summation by integration in (3.30).

The estimate which has just been obtained points to another interesting phenomenon. From the analysis of the two-level system summarized in (3.19) we know that a couple of levels can be considered as resonant, if the interaction strength $V$ between these levels exceeds the relative detuning $\Delta$. Naively
we could think that in the level-band system the same criterion is applicable, and therefore the number of the band states in resonance amounts to the number of levels $|n\rangle$ which falls into a strip of width $V$ around the level $|0\rangle$, that is $\mathcal{N}_{\text {naive }}=V / \delta$. However, the true result obtained reads $\mathcal{N}=(V / \delta)^{2}$. This implies that the states of the band do not interact with the level $|0\rangle$ independently but rather manifest a cooperative behavior: in order to obtain the correct estimate of the resonant levels with the help of the two-level condition $V>\Delta$ one has to replace the typical interaction $V$ by an effective (cooperative) interaction $V_{c}$ in the "naive" expression $\mathcal{N}=V_{c} / \delta$. This effective interaction $V_{c}$ equals the typical interaction $V$ multiplied by the square root of the true number of levels in resonance, that is $V \rightarrow V_{c}=V \sqrt{\mathcal{N}}$.

The origin of the $\sqrt{\mathcal{N}}$-factor is in the fact that consideration of the same problem from a different point of view in quantum mechanics usually implies a different choice of the basis sets, and all the transformations of the basis set are given by unitary matrices, which for all operators of physical values preserve sums of the absolute value squares of the matrix elements of any line and any column. Being applied to the operator of the dipole moment it yields the well-known optical theorem stating that the sum of squares of the transition dipole matrix elements that couple an arbitrary state to all other states is a constant, which remains the same even after reconstruction of the energy spectrum caused by a perturbation. This implies that the sum of squares of the transition matrix elements given by the product of a typical matrix element squared $V^{2}$ and the number $\mathcal{N}$ of the resonant states, that is $V^{2} \times(V / \delta)^{2}$, is a constant, and indeed the decay rate $W=\sqrt{V^{2} \times(V / \delta)^{2}}$ does not depend on the detailed structure of the band spectrum.

One can also look at the situation differently. Imagine that the levelband system originates from a two-level system, where one of the levels, say the level $|1\rangle$, interacts strongly with a $\mathcal{N}$-level band. As a result of this interaction each of the energy eigenstates $|n\rangle$ of the system "level $|1\rangle+\mathcal{N}$-level band" contains a contribution from the state $|1\rangle$, and the amplitude of this contribution $\langle 1 \mid n\rangle$ is of the order of $1 / \sqrt{\mathcal{N}}$ and hence matrix elements $V$ of the transitions from the state $|0\rangle$ to states $|n\rangle$ are $\sqrt{\mathcal{N}}$ times weaker, compared to the initial matrix element $V_{c}$ of the transition between the states $|0\rangle$ and $|1\rangle$. But it is clear that this matrix element $V_{c}$ has to be compared with a typical detuning $\Delta \simeq \delta \mathcal{N}$ in the two-level resonance criterion $V_{c} \sim \Delta \simeq \delta \mathcal{N}$, that is $V \sim \delta \sqrt{\mathcal{N}}$, or $\mathcal{N} \sim V^{2} / \delta^{2}$. Note that the decay rate $W=2 \pi V^{2} / \delta$ entering (3.34) can also be given in terms of the cooperative coupling $V_{c}=V \sqrt{\mathcal{N}}$ divided by a typical width $\Gamma=\delta \mathcal{N} / 2 \pi$ of the band, that is $W=V_{c}^{2} / \Gamma$.

## Wavepacket Representation

Sometimes it is very instructive to considerthe motion of the system shown in Fig. 3.4 in the coordinate representation, which allows one to see a deep analogy between the time evolution of the level-band system and the dynam-
ics of wavepackets. For this purpose one has to multiply the state amplitudes $\psi_{n}(t)$ by the coordinate-dependent eigenfunctions $\Psi_{n}(x)$ of the particle in the well $B$ and sum up over all $n$ in order to obtain the time-dependent wavefunction of the particle $\Psi(t, x)=\sum_{n} \psi_{n}(t) \Psi_{n}(x)$. After substitution of $\psi_{n}(t)=\frac{1}{2 \pi} \int_{-\infty+i \nu}^{+\infty+i \nu} e^{-i \varepsilon t} \psi_{n}(\varepsilon) d \varepsilon$ and $\psi_{n}(\varepsilon)=\psi_{0}(\varepsilon) V_{0}^{n} /\left(\varepsilon-\Delta_{n}\right)$ in accordance with (3.26) this yields

$$
\begin{equation*}
\Psi(t, x)=\frac{1}{2 \pi} \int_{-\infty+i \nu}^{+\infty+i \nu} e^{-i \varepsilon t} K(x, \varepsilon) \psi_{0}(\varepsilon) d \varepsilon \tag{3.41}
\end{equation*}
$$

where $K(x, \varepsilon)=\sum_{n} \Psi_{n}(x) V_{0}^{n} /\left(\varepsilon-\Delta_{n}\right)$. We note that the inverse Fourier transform of a product equals the convolution of the inverse Fourier transforms of the cofactors, and write this expression in the more convenient form

$$
\begin{equation*}
\Psi(t, x)=\int_{-\infty}^{t} K(x, t-\tau) \psi_{0}(\tau) d \tau \tag{3.42}
\end{equation*}
$$

where the time-dependent amplitude $\psi_{0}(t)$ of the state $|0\rangle$ is convoluted with the kernel $K(x, t)=\sum_{n} \Psi_{n}(x) V_{0}^{n} \exp \left\{-i \Delta_{n} t\right\}$, resulting from the inverse Fourier transformation of $K(x, \varepsilon)$.

When the spectrum of the band is dense, the total number of levels in the well is large, and the motion is one-dimensional, the semiclassical consideration provides us with a good expression for the wavefunction corresponding to a given energy $E_{n}$ and any smooth potential $U(x)$ :

$$
\begin{equation*}
\Psi_{n}(x)=\frac{1}{\sqrt[4]{E_{n}-U(x)}} \cos \left[\frac{\pi}{4}-\sqrt{2 m} \int_{0}^{x} \sqrt{E_{n}-U(z)} d z\right] \tag{3.43}
\end{equation*}
$$

where $x=0$ corresponds to the turning point, whereas for the rectangular potential of depth $\hbar \omega-U_{0}$ shown in Fig. 3.4 one has

$$
\begin{equation*}
\Psi_{n}(x)=\frac{1}{\sqrt[4]{\Delta_{n}+U_{0}}} \sin \left[\sqrt{2 m} x \sqrt{\Delta_{n}+U_{0}}\right] \tag{3.44}
\end{equation*}
$$

This yields the kernel

$$
\begin{equation*}
K(x, t)=\sum_{n} \frac{V_{0}^{n} \exp \left\{-i \Delta_{n} t\right\}}{\sqrt[4]{\Delta_{n}+U_{0}}} \sin \left[x \sqrt{2 m} \sqrt{\Delta_{n}+U_{0}}\right] \tag{3.45}
\end{equation*}
$$

which can be further simplified if we set as earlier all $V_{0}^{n}$ identical $V_{0}^{n}=V$, expand the root in the argument of sine in a Taylor series around $\Delta_{n}=0$, neglect indeed very slow dependence in the denominator and replace the sum by an integral by analogy to (3.31). We arrive at

$$
\begin{align*}
K(x, t) & =V \int \frac{\exp \{-i \delta t n\}}{\sqrt[4]{U_{0}}} \sin \left[\frac{\sqrt{m} x \delta}{\sqrt{2 U_{0}}} n\right] d n \\
& =\frac{i \pi V}{\delta \sqrt[4]{U_{0}}}[\delta(t-v x)-\delta(t+v x)] \tag{3.46}
\end{align*}
$$

where $v=\sqrt{2 U_{0}} / \sqrt{m}$ is the particle velocity in the right potential pit, and $\delta(t \pm v x)$ are Dirac $\delta$-functions. Since $x$ and $t$ are both positive only the first of these two $\delta$-functions is important. Substitution of (3.46) into (3.42) yields

$$
\begin{equation*}
\Psi(t, x)=\frac{i \pi V}{\delta \sqrt[4]{U_{0}}} \psi_{0}(t-v x) \tag{3.47}
\end{equation*}
$$

which means that the probability to detect the particle at the distance $x$ at the time moment $t$ equals (apart of a constant factor) the probability to detect the particle at the state $|0\rangle$ at the time $t-v x$, that is $\rho(t, x)=$ $|\Psi(t, x)|^{2} \sim \rho_{0}(t-v x)$.

## The Fano Problem

We now consider a more involved situation, when the continuum is not uniform, but is formed as a result of an interaction between a level and a uniform continuum shown in Fig. 3.7. It turns out that in such a situation the interference phenomena change significantly the Lorentzian line shapes of resonances. We consider the problem in the representation where the states of the band are not yet reduced to the eigenstate representation, but are still considered as a result of interaction of a state $|1\rangle$ with a uniform band of states $|k\rangle$.

The corresponding Schrödinger equation reads

$$
\begin{align*}
& i \dot{\psi}_{n}=\Delta_{n} \psi_{n}+V \psi_{0}+V^{\prime} \psi_{1} \\
& i \dot{\psi}_{0}=V \sum_{n=1}^{N} \psi_{n}+V^{\prime \prime} \psi_{1} \\
& i \dot{\psi}_{1}=V^{\prime} \sum_{n=1}^{N} \psi_{n}+\Delta_{1} \psi_{1}+V^{\prime \prime} \psi_{0} \tag{3.48}
\end{align*}
$$

where $\Delta_{1}$ denotes detuning between the states $|0\rangle$ and $|1\rangle$.
We now again proceed in the same way as earlier: after substitution of the Dirac $\delta$-function to the right-hand side of the second equation we perform the Fourier transformation and arrive at

$$
\begin{align*}
& \varepsilon \psi_{n}=\Delta_{n} \psi_{n}+V \psi_{0}+V^{\prime} \psi_{1} \\
& \varepsilon \psi_{0}=V \sum_{n=1}^{N} \psi_{n}+V^{\prime \prime} \psi_{1}+i \\
& \varepsilon \psi_{1}=V^{\prime} \sum_{n=1}^{N} \psi_{n}+\Delta_{1} \psi_{1}+V^{\prime \prime} \psi_{0} \tag{3.49}
\end{align*}
$$

We express $\psi_{n}$ in terms of $\psi_{0}$ and $\psi_{1}$ with the help of the first equation and find


Fig. 3.7. The Fano problem. (a) Representation of the band as a level interacting with a zone. The initial state $|0\rangle$ interacts with the state $|1\rangle$ by the coupling $V^{\prime \prime}$ and with the zone by the coupling $V$. The state $|1\rangle$ interacts with the zone by the coupling $V^{\prime}$. (b) Representation of the band in the energy eigenstates basis. The resulting coupling $V_{r}(E)$ is not uniform, but has the Fano profile $V_{r} \sim f(E)$ shown in (c).

$$
\begin{align*}
& \varepsilon \psi_{0}=V^{2} \sum_{n=1}^{N} \frac{1}{\left(\varepsilon-\Delta_{n}\right)} \psi_{0}+\left(V V^{\prime} \sum_{n=1}^{N} \frac{1}{\left(\varepsilon-\Delta_{n}\right)}+V^{\prime \prime}\right) \psi_{1}+i  \tag{3.50}\\
& \varepsilon \psi_{1}=\left(V^{\prime \prime}+V V^{\prime} \sum_{n=1}^{N} \frac{1}{\left(\varepsilon-\Delta_{n}\right)}\right) \psi_{0}+\left(\left(V^{\prime}\right)^{2} \sum_{n=1}^{N} \frac{1}{\left(\varepsilon-\Delta_{n}\right)}+\Delta_{1}\right) \psi_{1}
\end{align*}
$$

We now replace the sums by the integrals (3.50) and obtain

$$
\begin{align*}
\left(\varepsilon+i W_{0}\right) \psi_{0}-\left(V^{\prime \prime}-i W_{F}\right) \psi_{1} & =i \\
\left(V^{\prime \prime}-i W_{F}\right) \psi_{0}+\left(\Delta_{1}-i W_{1}-\varepsilon\right) \psi_{1} & =0 \tag{3.51}
\end{align*}
$$

where $W_{0}=\pi V^{2} / \delta ; W_{1}=\pi\left(V^{\prime}\right)^{2} / \delta ; W_{F}=\pi V V^{\prime} / \delta=\sqrt{W_{0} W_{1}}$. By solving (3.51) one obtains two expressions for $\psi_{0}$ and $\psi_{1}$ :

$$
\begin{align*}
& \psi_{0}(\varepsilon)=\frac{i\left(\Delta_{1}-i W_{1}-\varepsilon\right)}{\left(\varepsilon+i W_{0}\right)\left(\Delta_{1}-i W_{1}-\varepsilon\right)+\left(V^{\prime \prime}-i W_{F}\right)^{2}} \\
& \psi_{1}(\varepsilon)=\frac{-i\left(V^{\prime \prime}-i W_{F}\right)}{\left(\varepsilon+i W_{0}\right)\left(\Delta_{1}-i W_{1}-\varepsilon\right)+\left(V^{\prime \prime}-i W_{F}\right)^{2}} \tag{3.52}
\end{align*}
$$

which after substitution into the first equation of the set (3.49), yield

$$
\begin{equation*}
\psi_{n}=i \frac{\left(\Delta_{1}-i W_{1}-\varepsilon\right) V-\left(V^{\prime \prime}-i W_{F}\right) V^{\prime}}{\left(\varepsilon-\Delta_{n}\right)\left[\left(\varepsilon+i W_{0}\right)\left(\Delta_{1}-i W_{1}-\varepsilon\right)+\left(V^{\prime \prime}-i W_{F}\right)^{2}\right]} \tag{3.53}
\end{equation*}
$$



Fig. 3.8. The Fano problem. (a) Time-dependent populations of the initial state $\rho_{1}$ and the total population of the states $|0\rangle$ and $|1\rangle$. The Fano profile (b) of the population distribution over the states of the band.

As usual, the inverse Fourier transformation implies integration along the contour shown in Fig. 3.2, which can be carried out by finding the residuals in the singular points of the right-hand side of $(3.52),(3.53)$ at the lower part of the complex plane $\varepsilon$ including the real axis. The denominators of (3.52) are identical and have two singularities at the points $\varepsilon_{1}$ and $\varepsilon_{2}$, both of them with negative imaginary parts, that can be found as roots of a quadratic equation. Therefore each of the two amplitudes $\psi_{0}(t)$ and $\psi_{1}(t)$ manifests biexponential decay, whereas the corresponding populations $\rho_{0}(t)$ and $\rho_{1}(t)$ also
manifests an oscillatory interference pattern coming from the cross product of different singularities, as shown in Fig. 3.8(a). Each amplitude of the band states (3.53) also has an additional pole on the real axis; that at $t \rightarrow \infty$ does not decay but brings the main contribution of the form

$$
\begin{equation*}
\psi_{n}=i \frac{\left(\Delta_{1}-i W_{1}-\Delta_{n}\right) V-\left(V^{\prime \prime}-i W_{F}\right) V^{\prime}}{\left(\Delta_{n}+i W_{0}\right)\left(\Delta_{1}-i W_{1}-\Delta_{n}\right)+\left(V^{\prime \prime}-i W_{F}\right)^{2}} e^{i \Delta_{n} t} \tag{3.54}
\end{equation*}
$$

The corresponding population density distribution $f(\Delta)=\left|\psi_{n}\right|^{2} / \delta$ is no longer of the Lorentzian form (3.36), but has a typical two-hump shape known as the Fano profile, which is shown in Fig. 3.8(b). This profile also manifests itself in the dependence of the absorption rate on the frequency $\omega$ of a monochromatic excitation.

## Population Amplitude Injection

Let us now consider the level-band problem assuming that the states of this system were not populated initially, whereas the particle is pumped to the state $|0\rangle$ of the system in the course of an external process of a finite duration $\tau$. In other words the flux $\Pi(t)$ in (3.24) is not the Dirac $\delta$-function but a continuous function localized in a time interval of the order of $\tau$. The results of the consideration will show, on one hand, the relation between population distributions and line shapes, and from the other hand side will help us later on to understand the phenomena of population recurrences and revivals. Now instead of (3.24) we have

$$
\begin{align*}
i \dot{\psi}_{n} & =\Delta_{n} \psi_{n}+V_{n}^{0} \psi_{0} \\
i \dot{\psi}_{0} & =\sum_{n=1}^{N} V_{0}^{n} \psi_{n}+i \Pi(t) \tag{3.55}
\end{align*}
$$

which by analogy to (3.26) yields

$$
\begin{align*}
\psi_{0}(\varepsilon) & =\frac{i \Pi(\varepsilon)}{\varepsilon-\sum_{n=1}^{N} \frac{V_{0}^{n}}{\varepsilon-\Delta_{n}}} \\
\psi_{n}(\varepsilon) & =\frac{1}{\varepsilon-\Delta_{n}} V_{0}^{n} \psi_{0}=\frac{1}{\varepsilon-\Delta_{n}} \frac{i \Pi(\varepsilon) V_{0}^{n}}{\varepsilon-\sum_{n=1}^{N} \frac{V_{0}^{n}}{\varepsilon-\Delta_{n}}} . \tag{3.56}
\end{align*}
$$

Here $\Pi(\varepsilon)$ is the Fourier transform of $\Pi(t)$, which enters the equations for the amplitudes as an additional factor. We can now also identify the combination $V_{0}^{n} \psi_{0}$ with the probability amplitude supply corresponding to the transition from the state $|0\rangle$ to the state $|n\rangle$.

Let us now consider a particular case of the dependence $\Pi(t)$ by taking the Gaussian function

$$
\begin{equation*}
\Pi(t)=\frac{e^{-t^{2} / \tau^{2}}}{\sqrt{\pi} \tau} \tag{3.57}
\end{equation*}
$$

which yields

$$
\begin{equation*}
\Pi(\varepsilon)=\frac{1}{\sqrt{\pi} \tau} \int_{-\infty}^{\infty} e^{-t^{2} / \tau^{2}-i \varepsilon t} d t=e^{-\varepsilon^{2} \tau^{2} / 4} \tag{3.58}
\end{equation*}
$$

Substitution of (3.58) into (3.56) results, for the continuum,

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{i e^{-\varepsilon^{2} \tau^{2} / 4}}{\varepsilon-i V^{2} \pi / \delta} \tag{3.59}
\end{equation*}
$$



Fig. 3.9. A model of the population amplitude supply. (a) A particle with a broad wavepacket arrives from the left and penetrates to the state $|0\rangle$ through the left potential barrier with the penetration rate $\Pi(t)$ proportional to the probability amplitude to be near the barrier. The typical duration of the supply $\tau$ considerably exceeds the inverse decay rate $1 / W$ of the level. From the state $|0\rangle$ it comes to the states $|n\rangle$ of the right continuum (b). Interference of these states, taken with allowance of the coordinate dependence of their eigenfunctions, results in the propagation of a wavepacket $\Psi(x)$, which is so broad in space (and long in time), that the spectrum width (c) of the occupied states in the continuum is much less, compared to the Lorentzian width (3.36) (dotted line) given by the transparency of the second barrier $W$.

Let us consider the limiting case, when the decay rate exceeds the rate of particle supply, that is $V^{2} / \delta \gg 1 / \tau$. The Gaussian dependence makes $\psi_{0}(\varepsilon)$ vanishing for $\varepsilon>2 / \tau$ and therefore the term $\varepsilon$ in the denominator of (3.59) can be omitted, and we arrive at

$$
\begin{equation*}
\psi_{0}(\varepsilon)=-\frac{\delta}{V^{2} \pi} e^{-\varepsilon^{2} \tau^{2} / 4} \tag{3.60}
\end{equation*}
$$

which, after inverse Fourier transformation, reads

$$
\begin{equation*}
\psi_{0}(t)=-\frac{\delta}{V^{2} 2 \pi^{2}} \int_{\infty-i \nu}^{\infty+i \nu} e^{-\varepsilon^{2} \tau^{2} / 4+i \varepsilon t} d \varepsilon=-\frac{\delta}{V^{2} \pi \sqrt{\pi} \tau} e^{-t^{2} / \tau^{2}} \tag{3.61}
\end{equation*}
$$

One sees that in this case the amplitude of the probability to be at the level $|0\rangle$ follows the supply rate $\Pi(t)$, although it is scaled by the decay rate
factor $1 / W=\delta / V^{2} \pi$ and has opposite sign. This means that the probability amplitude of staying in the state $|0\rangle$ is attained as a result of balance of the supply process and the fast decay process corresponding to the rapid transmission of the probability amplitude to the states of the band, natural in the chosen limiting case. We illustrate this in Fig. 3.9 by an instructive analogy. We also note that $\Pi(\varepsilon)=$ const corresponds to the Dirac $\delta$-functionlike population amplitude injection, and $\Pi(\varepsilon) \sim 1 / \varepsilon$ allows for a steady state flux of the population amplitude.

### 3.2.3 Measurements and Relaxation as Processes in Level-Continuum Systems

The level-band system can serve as the simplest model of relaxation and of the measurement process in a two-level system. The compound system thus comprises the two-level system and either the meter or the environment. The latter has a very rich spectrum of eigenenergies, originating either from the classical character of the meter, or from the large number of degrees of freedom in the environment. States of the band are therefore related to the environment or the meter spectrum. The high band state density implies that the Heisenberg time of return is extremely long, and therefore both the measurement and the relaxation are irreversible processes, and we can consider the band as continuous.

We focus on the simplest case where the relaxation may occur only from the exited state $|1\rangle$ of the two-level system and that this system makes a transition to the ground state $|0\rangle$ as a result. This transition is accompanied by the transition in the environment associated with the energy transfer to one of many degrees of freedom. For a single elementary act of relaxation, we can say that the transition occurs from the vacuum state of the environment $|0,0,0, \ldots, 0\rangle$ to an excited state $|k\rangle=|0, \ldots, 1, \cdots, 0\rangle$ where the $k$-th degree of freedom gets excited from the ground to the first exited state. By $\psi_{0}$ we denote the probability amplitude for the compound system to be in the state $|0\rangle \otimes|0,0,0, \ldots, 0\rangle$, while $\psi_{1}$ and $\psi_{k}$ denote the amplitudes of the states $|1\rangle \otimes|0,0,0, \ldots, 0\rangle$ and $|0\rangle \otimes|k\rangle$, respectively.

Considering the measurement process, we assume that only the excited state of the two-level system can be detected. In other words, excitation of the two-level system provokes a transition in the meter, such that it leaves its vacuum state for an excited state $|k\rangle$, while the two-level system remains in the upper state $|1\rangle$. The complex variables $\psi_{0}, \psi_{1}$, and $\psi_{k}$ then denote the amplitudes of the states $|0\rangle \otimes|0\rangle,|1\rangle \otimes|0\rangle$, and $|1\rangle \otimes|k\rangle$, respectively.

Population dynamics is a coherent process, characterized by the probability amplitudes as far as the compound system is concerned. These amplitudes satisfy the Schrödinger equation

$$
\begin{align*}
i \dot{\psi}_{0} & =V \psi_{1} \\
i \dot{\psi}_{1} & =E_{1} \psi_{1}+V \psi_{0}+\sum_{k} V^{\prime} \psi_{k}  \tag{3.62}\\
i \dot{\psi}_{k} & =E_{k} \psi_{k}+V^{\prime *} \psi_{1}
\end{align*}
$$

where we have assumed that the upper and lower levels of the two-level system, separated by the energy $E_{1}$, are coupled by the interaction $V$, and that all the band states are coupled to the exited state by an identical interaction $V^{\prime}$. With the initial condition $\psi_{1}(t=0)=1$ after Fourier transformation this yields

$$
\begin{align*}
\psi_{0} & =\frac{i V}{\left(\varepsilon-E_{1}+i W(\varepsilon)\right) \varepsilon-V^{2}} \\
\psi_{1} & =\frac{i \varepsilon}{\left(\varepsilon-E_{1}+i W(\varepsilon)\right) \varepsilon-V^{2}}  \tag{3.63}\\
\psi_{k} & =\frac{V^{*}}{\varepsilon-E_{k}} \frac{i \varepsilon}{\left(\varepsilon-E_{1}+i W(\varepsilon)\right) \varepsilon-V^{2}}
\end{align*}
$$

where $W(\varepsilon)=i \sum_{k}\left|V^{\prime}\right|^{2}\left(\varepsilon-E_{k}\right)^{-1}$. For the continuum, infinitely large in both sides, the corresponding integration yields $W=\pi g\left|V^{\prime}\right|^{2}$ which does not depend on $\varepsilon$. After the inverse Fourier transformation we therefore arrive at

$$
\begin{align*}
\psi_{0}= & e^{-\left(i E_{1}+W\right) t / 2} \frac{2 V}{\Omega} \sin \frac{t \Omega}{2}  \tag{3.64}\\
\psi_{1}= & e^{-\left(i E_{1}+W\right) t / 2}\left[i \cos \frac{t \Omega}{2}+\frac{\left(E_{1}-i W\right)}{\Omega} \sin \frac{t \Omega}{2}\right] \\
\psi_{k}= & \frac{-i E_{k} V^{\prime *} e^{-i E_{k} t}}{E_{1} E_{k}+V^{2}-E_{k}\left(E_{k}+i W\right)}+\frac{V^{\prime *} e^{-\left(i E_{1}+W\right) t / 2}}{E_{1} E_{k}+V^{2}-E_{k}\left(E_{k}+i W\right)} \\
& \left.\times\left[i E_{k} \cos \Omega / 2+\left(E_{1} E_{k}+2 V^{2}-i E_{k} W\right)\right) \frac{\sin t \Omega / 2}{\Omega}\right]
\end{align*}
$$

where $\Omega=\sqrt{4 V^{2}+\left(E_{1}-i W\right)^{2}}$. Note that the absolute values of the amplitudes $\psi_{0}$ and $\psi_{1}$ decrease exponentially in time with the rate $(W-\operatorname{Im} \Omega) / 2$.

Equations (3.64) provide us with the solution for the compound system. However when concentrating only on the two-level system and ignoring the quantum structure of the environment or of the meter we need to take the trace of the compound system density matrix over all quantum numbers except those of the two levels considered. This yields

$$
\begin{align*}
\widehat{\rho}_{r} & =\left(\psi_{0}|0\rangle+\psi_{1}|1\rangle\right)\left(\psi_{0}^{*}\langle 0|+\psi_{1}^{*}\langle 1|\right)+\sum \psi_{k}^{*} \psi_{k}|0\rangle\langle 0| \\
\hat{\rho}_{m} & =\left(\psi_{0}|0\rangle+\psi_{1}|1\rangle\right)\left(\psi_{0}^{*}\langle 0|+\psi_{1}^{*}\langle 1|\right)+\sum \psi_{k}^{*} \psi_{k}|1\rangle\langle 1|, \tag{3.65}
\end{align*}
$$

for the environment and the meter, respectively. When we take into account the normalization condition $\psi_{0}^{*} \psi_{0}+\psi_{1}^{*} \psi_{1}+\sum \psi_{k}^{*} \psi_{k}=1$ for the compound system, (3.65) adopts the form

$$
\begin{align*}
\hat{\rho}_{r} & =\psi_{0}^{*} \psi_{1}|1\rangle\langle 0|+\psi_{1}^{*} \psi_{0}|0\rangle\langle 1|+\psi_{1}^{*} \psi_{1}|1\rangle\langle 1|+\left(1-\psi_{1}^{*} \psi_{1}\right)|0\rangle\langle 0| \\
\widehat{\rho}_{m} & =\psi_{0}^{*} \psi_{1}|1\rangle\langle 0|+\psi_{1}^{*} \psi_{0}|0\rangle\langle 1|+\psi_{0}^{*} \psi_{0}|0\rangle\langle 0|+\left(1-\psi_{0}^{*} \psi_{0}\right)|1\rangle\langle 1| . \tag{3.66}
\end{align*}
$$

One sees that, as a result of the elementary act of relaxation, the two-level system comes to the state $|0\rangle$ and the measured system comes to the state $|1\rangle$ after a time $1 /(W-\operatorname{Im} \Omega)$, when both amplitudes $\psi_{0}$ and $\psi_{1}$ vanish.


Fig. 3.10. Dynamics of the Bloch vector $B=\left(\rho_{01}+\rho_{10}, i \rho_{10}-i \rho_{01}, \rho_{11}-\rho_{00}\right)$ representing the density matrix. (a) For a relaxing two-level system the vector $B$ starts at the point $(1,0,0)$ corresponding to the excited state and goes to the point $(-1,0,0)$ following the spiral trajectory associated with Rabi oscillations. For a system measured in the upper states the vector returns to the initial position along a spiral trajectory. (b) Exponential relaxation acquires oscillatory behavior in the case of a finite bandwidth. (c) The finite bandwidth also affects the trajectory of $B$.

A Bloch vector $B=\left(2 \operatorname{Re} \rho_{10}, 2 \operatorname{Im} \rho_{10}, \rho_{11}-\rho_{00}\right)$ composed of the density matrix elements (3.66) with the allowance of amplitudes (3.64) is a convinient real quantity which helps to trace the time evolution of an open two-level system. Trajectories of $B$ shown in Fig. 3.10a) illustrate the dynamics of the two-level system, initially in the upper state $|1\rangle$, which is subjected to either measurement or relaxation. One sees that the vector, initially in the utmost upper position $\rho_{11}=1$, starts to move toward lower positions as a result of the interaction $V$ but returns along a spiral trajectory as a result of the elementary act of measurement. As a result of the elementary act of relaxation the vector gradually moves to the lowest position $\rho_{00}=1$, also along a spiral trajectory.

Let us derive differential equations for the density matrices $\widehat{\rho}_{r}$ and $\widehat{\rho}_{m}$. We note that in the case of an infinite band one finds

$$
\begin{align*}
i \dot{\psi}_{0} & =V \psi_{1}  \tag{3.67}\\
i \dot{\psi}_{1} & =E_{1} \psi_{1}+V \psi_{0}-i W \psi_{1}
\end{align*}
$$

which results from the Fourier transformation of the last two equations (3.62), substitution of $\psi_{k}(\varepsilon)$ from the third equation to the second one, and the replacement $\sum_{k} \frac{\left|V^{\prime}\right|^{2}}{\varepsilon-E_{k}} \rightarrow-i W$ followed by the inverse Fourier transformation. With the help of (3.67) and their complex conjugate, after direct calculations one obtains the differential equations

$$
\begin{gather*}
-i \frac{\partial}{\partial t}\left(\begin{array}{cc}
1-\psi_{1}^{*} \psi_{1} \psi_{0}^{*} \psi_{1} \\
\psi_{1}^{*} \psi_{0} & \psi_{1}^{*} \psi_{1}
\end{array}\right)= \\
\left(\begin{array}{cc}
V \psi_{1}^{*} \psi_{0}-V \psi_{0}^{*} \psi_{1}-i 2 W \psi_{1}^{*} \psi_{1} & \left(i W-E_{1}\right) \psi_{0}^{*} \psi_{1}-V\left(\psi_{0}^{*} \psi_{0}-\psi_{1}^{*} \psi_{1}\right) \\
\left(i W+E_{1}\right) \psi_{1}^{*} \psi_{0}-V\left(\psi_{1}^{*} \psi_{1}-\psi_{0}^{*} \psi_{0}\right) & V \psi_{0}^{*} \psi_{1}-V \psi_{1}^{*} \psi_{0}+i 2 W \psi_{1}^{*} \psi_{1}
\end{array}\right) \\
-i \frac{\partial}{\partial t}\left(\begin{array}{cc}
\psi_{0}^{*} \psi_{0} & \psi_{0}^{*} \psi_{1} \\
\psi_{1}^{*} \psi_{0} & 1-\psi_{0}^{*} \psi_{0}
\end{array}\right)= \\
\left(\begin{array}{cc}
V \psi_{1}^{*} \psi_{0}-V \psi_{0}^{*} \psi_{1} & \left(i W-E_{1}\right) \psi_{0}^{*} \psi_{1}-V\left(\psi_{0}^{*} \psi_{0}-\psi_{1}^{*} \psi_{1}\right) \\
\left(E_{1}+i W\right) \psi_{1}^{*} \psi_{0}-V\left(\psi_{1}^{*} \psi_{1}-\psi_{0}^{*} \psi_{0}\right) & V \psi_{0}^{*} \psi_{1}-V \psi_{1}^{*} \psi_{0}
\end{array}\right) \tag{3.68}
\end{gather*}
$$

which contain all four binary combinations of the amplitudes on the righthand side. Therefore it is impossible to express the right-hand side only in terms of the matrix elements of $\widehat{\rho}_{r}$ or $\widehat{\rho}_{m}$ each of which contains only three such combinations. However one can write these equations in the form

$$
\begin{align*}
& i \frac{\partial}{\partial t} \widehat{\rho}_{r}=\left[H, \widehat{\rho}_{r}\right]-i W\left(\begin{array}{cc}
-2 \rho_{r 11} & \rho_{r 10} \\
\rho_{r 01} & 2 \rho_{r 11}
\end{array}\right)-V\left(\begin{array}{cc}
0 & 1 \\
-1 & 0
\end{array}\right) \sum \psi_{k}^{*} \psi_{k}  \tag{3.69}\\
& i \frac{\partial}{\partial t} \widehat{\rho}_{m}=\left[H, \widehat{\rho}_{m}\right]-i W\left(\begin{array}{cc}
0 & \rho_{m 10} \\
\rho_{m 01} & 0
\end{array}\right)+V\left(\begin{array}{cc}
0 & 1 \\
-1 & 0
\end{array}\right) \sum \psi_{k}^{*} \psi_{k}
\end{align*}
$$

involving the sum over the band states.
In the situation where many elementary acts of relaxation or measurements can occur continuously, one can omit the terms $\psi_{k}^{*} \psi_{k}$ on the righthand side of both equations (3.69), since immediately after the transition to a band state $|k\rangle$ the system starts to be subjected to the influence of the next meter or the next set of degrees of freedom of the environment in the vacuum state, as if it were transported back to the states corresponding to the two-level system and the vacuum state of the previous meter or environment. This gives the standard form known as the master equations. However, for the consideration of a single elementary act of measurement or relaxation the sum on the right-hand side of (3.69) is important. In order to express it in terms of the density matrix elements we make use of the normalization condition and take the time derivative of the sum

$$
\begin{gather*}
\sum \psi_{k}^{*} \psi_{k}=1-\psi_{0}^{*} \psi_{0}-\psi_{1}^{*} \psi_{1} \\
i \frac{\partial}{\partial t} \sum \psi_{k}^{*} \psi_{k}=-i \frac{\partial}{\partial t}\left(\psi_{0}^{*} \psi_{0}+\psi_{1}^{*} \psi_{1}\right)=-i 2 W \psi_{1}^{*} \psi_{1} \tag{3.70}
\end{gather*}
$$

which allows one to eliminate the unknown functions and write the second order differential equations

$$
\begin{align*}
- & \frac{\partial^{2}}{\partial t^{2}} \widehat{\rho}_{r}=\left[H, i \frac{\partial}{\partial t} \widehat{\rho}_{r}\right]-W \frac{\partial}{\partial t}\left(\begin{array}{cc}
2 \rho_{r 11} & -\rho_{r 10} \\
-\rho_{r 01} & -2 \rho_{r 11}
\end{array}\right)+i 2 W V\left(\begin{array}{cc}
0 & \rho_{r 11} \\
-\rho_{r 11} & 0
\end{array}\right), \\
& \left(\frac{\partial}{\partial t}-2 W\right) \frac{\partial}{\partial t} \widehat{\rho}_{m}=-i\left[H,\left(\frac{\partial}{\partial t}-2 W\right) \widehat{\rho}_{m}\right] \\
- & W\left(\frac{\partial}{\partial t}-2 W\right)\left(\begin{array}{cc}
0 & \rho_{m 10} \\
\rho_{m 01} & 0
\end{array}\right)+i 2 V W\left(\begin{array}{cc}
0 & \rho_{m 11} \\
-\rho_{m 11} & 0
\end{array}\right) \tag{3.71}
\end{align*}
$$

for the density matrices $\widehat{\rho}_{r}$ and $\widehat{\rho}_{m}$ (3.66). These equations have to be employed in the situation where the environment or the meter saturates after a single elementary act of relaxation or measurement, respectively.

Note the important particular regime of rapid and continuous measurements. After having neglected the term $\psi_{k}^{*} \psi_{k},(3.69)$ for $\hat{\rho}_{m}$ takes the form

$$
\begin{align*}
& i \frac{\partial}{\partial t}\left(\begin{array}{cc}
\rho_{00} & \rho_{m 10} \\
\rho_{m 01} & \rho_{11}
\end{array}\right)=  \tag{3.72}\\
& \left(\begin{array}{cc}
V \rho_{10}-V \rho_{01} & -\left(i W-E_{1}\right) \rho_{10}+V\left(\rho_{00}-\rho_{11}\right) \\
-\left(E_{1}+i W\right) \rho_{01}-V\left(\rho_{00}-\rho_{11}\right) & -V \rho_{m 10}+V \rho_{m 01}
\end{array}\right)
\end{align*}
$$

and in the case $W \gg V$ this yields

$$
\begin{gather*}
\rho_{01}=V \frac{\rho_{11}-\rho_{00}}{E_{1}+i W} ; \quad \rho_{10}=V \frac{\rho_{00}-\rho_{11}}{i W-E_{1}} \\
\frac{\partial}{\partial t}\left(\rho_{11}-\rho_{00}\right)=\frac{-4 W V^{2}}{E_{1}^{2}+W^{2}}\left(\rho_{11}-\rho_{00}\right) \tag{3.73}
\end{gather*}
$$

which means that the off-diagonal elements of the density matrix rapidly assume their stationary values given by the instantaneous value of the difference $\rho_{11}-\rho_{00}$ while the latter change slowly. One sees that for $W \gg E_{1}$ the relaxation rate $4 V^{2} / W$ of the diagonal matrix elements is inversely proportional to the measurement rate $W$ and tends to zero when the measurements are performed infinitely fast. This implies that the system originally in the state $|1\rangle$ remains in this state indefinitely. This phenomenon is known as the Zeno effect and can be observed in any quantum system subjected to frequent measurements, where the relaxation rate of the measured state is inversely proportional to the rate of measurement.

Now let us consider the role of the spectral characteristics of the environment and the meter, which have been taken as an infinite continuum in the previous analysis. If the latter is not the case, the relaxation rate $W$ is not a constant, but a function which depends on the frequency $\varepsilon$. Therefore in the time-dependent form of the Schrödinger equation it corresponds to an integral operator of convolution, with the kernel given by the Fourier transform of $W(\varepsilon)$. In the simplest case $W(\varepsilon)=w /(\varepsilon-i \Gamma)$, soluble analytically, one finds

$$
\begin{align*}
\psi_{0} & =i \frac{V}{\varepsilon\left(\varepsilon-E_{1}-w /(\varepsilon+i \Gamma)\right)-V^{2}} \\
\psi_{1} & =\frac{i \varepsilon}{\left(\varepsilon-E_{1}-w /(\varepsilon+i \Gamma)\right) \varepsilon-V^{2}} \tag{3.74}
\end{align*}
$$

These expressions have three simple poles each, in contrast to (3.63) where the same quantities had only two poles. The additional pole and the associated frequency account for the dynamics of the meter or environment resulting from the non-instantaneous character of their motion. The narrower the band width $\Gamma$, the larger the influence of this motion on the process of relaxation and measurement. In Fig.3.10(b) we illustrate the difference between relaxation to an infinite band and to a band of a finite width for the example of the population difference between the upper and lower levels suggested by the Fourier transforms of (3.74) and the definitions (3.66). The time dependence of the population difference is exponential for the case of an infinite band, but it experiences dying oscillations when the band is narrow. Other elements of the density matrix represented in Fig.3.10(c) in the form of the Bloch vector $B$ are also affected by the band width.

### 3.3 Long-Time Behavior

We are now in the position to consider dynamics of the level-band system (3.23) at the long time limit $t \gg 1 / \delta$, that is at time longer compared to the Heisenberg return time from the band, or speaking in terms of the example shown in Fig. 3.4, longer than the period $T$ of oscillations of the particle in the potential pit $B$. It turns out that in a certain regime one observes at this timescale a number of interesting universal phenomena: quantum recurrences, revivals, and fractional revivals. We start consideration of the long-time limit with a general formalism, and later we turn to the phenomenon of quantum recurrences, which occurs first, that is at a time-scale, shorter than that for two other phenomena - revivals and fractional revivals.

### 3.3.1 General Consideration of the Long-Time Limit

For the sake of simplicity we restrict ourselves to the case when the population amplitude has been injected to the system during a time $\tau$, which is shorter than the period $T \sim 1 / \delta$, but still much longer than the typical decay time $1 / W \sim \delta / V^{2}$ of the state $|0\rangle$. The reason for this is purely technical, since in this case the analytical expressions for the amplitudes are much simpler than for the general case, but nevertheless they manifest all physically important phenomena typical of the level-band system. Full analysis, which can be done by analogy, results in much more cumbersome expressions. For the same reason we restrict ourselves only to the case of identical transition amplitudes $V_{n}^{0}=V_{0}^{n}=V$, and an equidistant spectrum of the band $\Delta_{n}=\delta n$, where $n$ takes all integer values, positive and negative. In such a case one of the states of the band, $n=0$, has the same energy and the same notation as the level $|0\rangle$, which hopefully will not create any confusion, since this state always enters expressions in combination with other states $|n\rangle$.

Let us concentrate on zeros of the denominator of the first of (3.26), that have to be found from the solution of the algebraic equation which now reads

$$
\begin{equation*}
\varepsilon-\sum_{n=-\infty}^{\infty} \frac{V^{2}}{\varepsilon-\delta n}=0 \tag{3.75}
\end{equation*}
$$

and can also be written in the form

$$
\begin{equation*}
\frac{\varepsilon}{\delta}=\frac{\pi V^{2}}{\delta^{2}} \cot \left(\pi \frac{\varepsilon}{\delta}\right) \tag{3.76}
\end{equation*}
$$

since $\sum_{n=-\infty}^{\infty}(x-n \pi)^{-1}=\cot x$. In Fig. 3.11 we depict the right- and the left-hand sides of (3.76), where the roots correspond to the points of crossing. For large $V^{2} / \delta^{2}$ the crossing occurs near the points $\varepsilon_{n} \simeq(n+1 / 2) \delta$, where


Fig. 3.11. Energy eigenstates of the level-band system corresponding to crossing points of the linear function $\varepsilon$ and the band resolvent $\sum_{n} V^{2} /\left(\varepsilon-\Delta_{n}\right)=$ $\left(\pi V^{2} / \delta\right) \cot (\pi \varepsilon / \delta)$.
$\cot (\varepsilon \pi / \delta)$ is small. Let us denote by $\alpha_{n} \delta / \pi$ the deviation of the $n$-th root $\varepsilon_{n}$ from the point $(n+1 / 2) \delta$, and write (3.76) in the form

$$
\begin{align*}
\left(n+\frac{1}{2}\right)+\frac{\alpha_{n}}{\pi} & =\frac{\pi V^{2}}{\delta^{2}} \cot \left(\pi\left(n+\frac{1}{2}\right)+\alpha_{n}\right) \\
& =-\frac{\pi V^{2}}{\delta^{2}} \tan \left(\alpha_{n}\right) \simeq-\frac{\pi V^{2}}{\delta^{2}} \alpha_{n}\left(1+\frac{\alpha_{n}^{2}}{3}\right) \tag{3.77}
\end{align*}
$$

We consider $\alpha_{n} \sim \arctan \left(\delta^{2} n / \pi V^{2}\right)$ as small, since the typical number $n \sim$ $1 / \tau \delta$ corresponding to a band state populated in the course of the spectrally localized $(\sim 1 / \tau)$ injection remains small as compared to the total number of states $\pi V^{2} / \delta^{2}$ mixed up by the interaction $V$, as shown in Fig. 3.9.

The third-order equation (3.77) can be resolved exactly with respect to $\alpha_{n}$. However, to be consistent with the order of approximation we have to write it also in the approximative form

$$
\begin{equation*}
\alpha_{n}=\frac{\pi^{5} V^{2}}{\delta^{2}} \frac{\left(n+\frac{1}{2}\right)^{3}}{3\left(1+\frac{\pi^{2} V^{2}}{\delta^{2}}\right)^{4}}-\frac{\pi\left(n+\frac{1}{2}\right)}{1+\frac{\pi^{2} V^{2}}{\delta^{2}}}, \tag{3.78}
\end{equation*}
$$

which for the same order of approximation yields

$$
\begin{equation*}
\varepsilon_{n}=\left(n+\frac{1}{2}\right) \delta^{\prime}+\alpha\left(n+\frac{1}{2}\right)^{3} \tag{3.79}
\end{equation*}
$$

where $\delta^{\prime}=\delta\left(1-\delta^{2} / \pi^{2} V^{2}\right)$ and $\alpha=\delta^{7} / 3 \pi^{4} V^{6}$.
What is the physical meaning of (3.79)? It gives the position of energy eigenstates of the system of the level and the band and implies that the density of states $1 / \delta^{\prime}$ of such a system is slightly larger than the state density $1 / \delta$ of the band, since the same spectral domain now contains one more level. Moreover, it shows that the energy eigenstates are not equidistant, as was the case for the band, but experience a third-order correction $\alpha(n+1 / 2)^{3}$. One can also notice than $n$ enters the equation in the combination $(n+1 / 2)$, which implies that for the strong perturbation the energy eigenstates are located in the middle of the intervals among the successive band states.

In order to find the time evolution of the system we have to perform the inverse Fourier transformation, which implies calculation of the residuals at the points of the singularities. One needs to calculate derivatives of the denominator (3.75) for this purpose. We note that $d(\cot y x) / d x=-y / \sin ^{2} y x$, calculate the derivative and obtain

$$
\begin{equation*}
1-\frac{\pi V^{2}}{\delta} \frac{\pi}{\delta} \frac{1}{\sin ^{2} \frac{\pi \varepsilon_{n}}{\delta}} \simeq-\frac{\pi^{2} V^{2}}{\delta^{2}} \tag{3.80}
\end{equation*}
$$

where we have substituted (3.79) and took into account the relation $\tau V^{2} / \delta \gg$ 1 that is $V^{2} / \delta^{2} \gg n$. We now make use of the formula

$$
\begin{equation*}
\frac{1}{f(x)}=\sum_{x_{n}} \frac{1}{\left(x-x_{n}\right) f^{\prime}\left(x_{n}\right)} \tag{3.81}
\end{equation*}
$$

which is valid for any analytical function $f(x)$ that has only simple poles in the points $x_{n}$, and which implies summation over all these points. This formula allows one to write (3.56)

$$
\begin{equation*}
\psi_{0}(\varepsilon)=-\frac{i \delta^{2} \Pi(\varepsilon)}{\pi^{2} V^{2}} \sum_{n} \frac{1}{\varepsilon-\left(n+\frac{1}{2}\right) \delta^{\prime}-\alpha\left(n+\frac{1}{2}\right)^{3}} \tag{3.82}
\end{equation*}
$$

where (3.79) and (3.80) have also been employed. We also give here without derivation the expression

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\sum_{n} \frac{-i \delta \Pi(\varepsilon) V^{2}}{\left(\pi V^{2} / \delta\right)^{2}+\left[\left(n+\frac{1}{2}\right) \delta^{\prime}\right]^{2}} \frac{1}{\left[\varepsilon-\left(n+\frac{1}{2}\right) \delta^{\prime}-\alpha\left(n+\frac{1}{2}\right)^{3}\right]} \tag{3.83}
\end{equation*}
$$

which is also valid in the case $\tau V^{2} / \delta<1$ and relies only on the requirement $V^{2} / \delta^{2} \gg 1$.

Note that expressions (3.81) and (3.82) have a clear quantum mechanical meaning for the case $\Pi(t)=\delta(t)$, that is $\Pi(\varepsilon)=1$ : we just project the timedependent quantum state of the system given in terms of the eigenstates $\left|\varepsilon_{n}\right\rangle$ to the initial state $|0\rangle$, which after Fourier transformation yields the amplitude

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\sum_{n} \frac{\left\langle 0 \mid \varepsilon_{n}\right\rangle}{\varepsilon-\varepsilon_{n}} \tag{3.84}
\end{equation*}
$$

whereas the projections $\left\langle 0 \mid \varepsilon_{n}\right\rangle$ are given in terms of the residuals $1 / f_{\varepsilon}^{\prime}(\varepsilon)$ at the points $\varepsilon=\varepsilon_{n}$.

Inverse Fourier transformation of (3.82) recovers the time-dependent amplitude $\psi_{0}(t)$ and yields

$$
\begin{equation*}
\psi_{0}(t)=\frac{2 \delta^{2}}{\pi V^{2}} \sum_{n} e^{-\left[(n+1 / 2) \delta^{\prime}-\alpha(n+1 / 2)^{3}\right]^{2} \tau^{2} / 4-i(n+1 / 2) \delta^{\prime} t-i \alpha(n+1 / 2)^{3} t} \tag{3.85}
\end{equation*}
$$

for the case of $\Pi(t)$ given by (3.57). We note that for a typical size of $n \sim 1 / \tau \delta$ and for $\alpha \sim \delta^{7} / V^{6}$ the combination $\alpha\left(n+\frac{1}{2}\right)^{3} \sim\left(\delta / V^{2} \tau\right)^{3} \tau \delta$ can be omitted as a product of two small parameters, and we are left with

$$
\begin{equation*}
\psi_{0}(t)=\frac{\delta^{2}}{\pi^{2} V^{2}} \sum_{n} e^{-(n+1 / 2)^{2} \delta^{\prime 2} \tau^{2} / 4-i(n+1 / 2) \delta^{\prime} t-i \alpha(n+1 / 2)^{3} t} \tag{3.86}
\end{equation*}
$$

If we decide to consider the case $t \ll 1 / \delta$, neglect the cubic term in the exponent and replace summation by integration, we immediately recover the result (3.61) which has been obtained directly under the same assumption. Our aim here is to go beyond the limit $t \ll 1 / \delta$, which does not allow one to replace sums by integrals, since the phase difference between the neighboring terms of the sum (3.86) is of the order of $\pi$.

### 3.3.2 Quantum Recurrences

Generally speaking all three terms in the exponent (3.86) are equally important, since they depend on different combinations of parameters. Nevertheless one can identify a regime $t \sim 1 / \delta^{\prime}$, where the cubic term is still negligible, since $n \sim 1 / \tau \delta$, and hence $\alpha\left(n+\frac{1}{2}\right)^{3} t \sim\left(\delta / V^{2} \tau\right)^{3} t \delta \sim\left(\delta / V^{2} \tau\right)^{3} \ll 1$. In terms of the wavepackets shown in Figs. 3.4 and 3.9 this regime corresponds to the situation, where the wavepacket emitted by the level starts to return after either one or several oscillation periods. We therefore have the expression

$$
\begin{equation*}
\psi_{0}(t)=\frac{\delta^{2}}{\pi^{2} V^{2}} \sum_{n} e^{-(n+1 / 2)^{2} \delta^{\prime 2} \tau^{2} / 4-i(n+1 / 2) \delta^{\prime} t} \tag{3.87}
\end{equation*}
$$

where summation cannot be replaced by integration.
How do we calculate this sum? Consideration of wavepackets give us an idea about the structure of the sum: it consists of a series of spikes corresponding to the time moments when the wavepacket returns to the left turning point after each oscillation period. A mathematical tool known as the Poisson summation formula turns out to be very convenient for calculation of the sum in such a case. Apparently one can write (3.87) in the form of an integral

$$
\begin{equation*}
\psi_{0}(t)=\frac{\delta^{2}}{\pi^{2} V^{2}} \int_{-\infty}^{\infty} \sum_{n} \delta(x-n) e^{-(x+1 / 2)^{2} \delta^{\prime 2} \tau^{2} / 4-i(x+1 / 2) \delta^{\prime} t} d x \tag{3.88}
\end{equation*}
$$

since the sum $\sum_{n} \delta(x-n)$ of Dirac $\delta$-functions just pickups the values of the integrand in the integer points. Now we note that

$$
\begin{equation*}
\sum_{n} \delta(x-n)=\sum_{k} e^{i 2 \pi k x} \tag{3.89}
\end{equation*}
$$

which can be proved just by casting both sides of this equation, which is periodic in $x$, into Fourier series on any interval of length 1 placed around any integer point. Indeed, for such an interval which starts at a point $O$ we find the $m$-th harmonic of the left-hand side

$$
\begin{equation*}
\int_{O}^{O+1} e^{i 2 \pi m x} \sum_{n} \delta(x-n) d x=\int_{O}^{O+1} e^{i 2 \pi m x} \delta\left(x-n_{O}\right) d x=e^{i 2 \pi m n_{O}}=1 \tag{3.90}
\end{equation*}
$$

and of the right-hand side

$$
\begin{equation*}
\int_{O}^{O+1} e^{i 2 \pi m x} \sum_{k} e^{i 2 \pi k x} d x=\sum_{k} \int_{O}^{O+1} e^{i 2 \pi(m-n) x} d x=\sum_{k} \delta_{n}^{m}=1 \tag{3.91}
\end{equation*}
$$

where $n_{O}$ is the integer part of $O$, and $\delta_{n}^{m}$ is the Kronecker symbol. Substitution of (3.89) into (3.88) yields

$$
\begin{equation*}
\psi_{0}(t)=\frac{\delta^{2}}{\pi^{2} V^{2}} \sum_{k} \int_{-\infty}^{\infty} e^{-(x+1 / 2)^{2} \delta^{\prime 2} \tau^{2} / 4-i(x+1 / 2) \delta^{\prime} t+i 2 \pi k x} d x \tag{3.92}
\end{equation*}
$$

where we have changed the order of the summation and the integration. The Gaussian integration is straightforward and with the allowance of the integral (3.61) for $\varepsilon=(x+1 / 2) \delta^{\prime}$ results in the sum

$$
\begin{equation*}
\psi_{0}(t)=\frac{\delta}{2 \pi V^{2}} \sum_{k} \frac{e^{-i \pi k}}{\sqrt{\pi} \tau} \exp \left\{-\frac{\left(t-2 \pi k / \delta^{\prime}\right)^{2}}{\tau^{2}}\right\} \tag{3.93}
\end{equation*}
$$

where we have neglected the difference between $\delta^{\prime}$ and $\delta$ in the pre-exponential factor.

One recognizes in (3.93) a series of Gaussian spikes, which apart of the phase factor $e^{i \pi k}=(-1)^{k}$ differ only by their arrival times $2 \pi k / \delta^{\prime}$ equally spaced in time by the interval $T_{r}=2 \pi / \delta^{\prime}$, which for the model in Fig. 3.4 has an evident meaning of the oscillation period for the particle in the well $B$. We illustrate the situation in Fig. 3.12. The wavepacket created as the result of Gaussian injection of the population amplitude to the level $|0\rangle$ propagates toward the right turning point. It reaches the point at $t \simeq \pi / \delta$ and after reflection with phase inversion starts to propagate back. At $t \simeq 2 \pi / \delta$ it returns and penetrates through the barrier with the extent of penetration given by the ratio $\sim W^{-1} / \tau$ of the decay time $W^{-1}=\delta / \pi V^{2}$ and the injection time $\tau$.


Fig. 3.12. Recurrences in the level-band system corresponding to a two-well potential. (a) Injection of the population amplitude to the system and creation of the packet. (b) After half of the oscillation period, reflection of the wavepacket with $\pi$-shift of the phase. (c) The recurrence. The presence of the particle at the initial level is restricted by the necessity of penetration through the barrier.

In the extreme case we have just considered the wavepacket of the returned particle has precisely the same shape as it is had at the beginning of the evolution. This is not always the case in the presence of anharmonic (that is nonlinear in $n$ ) terms that have been omitted in our consideration. The $n^{2}$ term results in dispersion of the wavepacket, which manifests itself in spreading, that is in the increase of the recurrence duration. It has not been taken into account initially in our band model with equidistant levels. Although the $n^{3}$ anharmonicity of the band has also not been taken into account, it arises as a result of the interaction of the band states with the level $|0\rangle$ in the form $\alpha(n+1 / 2)^{3}$, as one can see in (3.86), and results in more significant transformation of the shape of spikes at long time, which is going to be considered now.

The third-order anharmonicity plays the dominating role for the time interval $1 / \alpha \gg t \gg 1 / \alpha n^{3} \sim 1 / \alpha^{\prime}(\tau \delta)^{3}$, that is $(V / \delta)^{6} \gg \delta t \gg\left(V^{2} / \delta\right) \tau$, where one can neglect the square term compared to the cubic term in the exponent of (3.86). The Poisson summation formula for this case yields

$$
\begin{equation*}
\psi_{0}(t)=\frac{\delta^{2}}{\pi^{2} V^{2}} \sum_{k} e^{-i \pi k} \int_{-\infty}^{\infty} e^{-i \alpha(x+1 / 2)^{3} t-i(x+1 / 2) \delta^{\prime} t+i 2 \pi k(x+1 / 2)} d x \tag{3.94}
\end{equation*}
$$

where after replacement $x+1 / 2 \rightarrow x$ and regrouping of terms linear in $x$ in the exponent one recognizes in the integral on the right-hand side an integral representation of the Airy function $2 \sqrt{\pi} \operatorname{Ai}(z)=\int_{-\infty}^{\infty} \exp \left(i x^{3} / 3+i x z\right) d x$ and arrives at

$$
\begin{equation*}
\psi_{0}(t)=\frac{\delta^{2}}{\pi \sqrt{\pi} V^{2} \sqrt[3]{3 \alpha t}} \sum_{k}(-1)^{k} \operatorname{Ai}\left(\frac{2 \pi k-\delta^{\prime} t}{\sqrt[3]{3 \alpha t}}\right) \tag{3.95}
\end{equation*}
$$



Fig. 3.13. A series of population recurrences for $V^{2} / \delta \simeq 1.5$ and $\tau \delta \simeq 0.6$ with corresponding anharmonicity $\alpha=0.001$. The time axis is expressed in the numbers of periods. One clearly sees a cross-over from the bell-shape spikes to the shapes with asymmetric oscillations, typical of the Airy function.

Actually, this summation can be performed for the exact expression (3.86), and even in a more general case, when along with the cubic nonlinearity $\alpha(n+1 / 2)^{3} t$ a quadric term $\beta(n+1 / 2)^{2} t$ enters the phase of the integrand. One gets rid of the quadric term just by a shift of the integration variable $x \rightarrow x-i a$. The final expression reads

$$
\begin{equation*}
\psi_{0}(t)=\frac{\delta^{2} e^{2 a^{3} / 3}}{\pi \sqrt{\pi} V^{2} \sqrt[3]{3 \alpha t}} \sum_{k}(-1)^{k} e^{a b} \operatorname{Ai}\left(b+a^{2}\right) \tag{3.96}
\end{equation*}
$$

where $a=\left(-\delta^{\prime 2} \tau^{2} / 4+i \beta t\right) /(\sqrt[3]{3 \alpha t})^{2}$ and $b=\left(2 \pi k-\delta^{\prime} t\right) / \sqrt[3]{3 \alpha t}$ enter the argument of the Airy function in (3.95).

In Fig. 3.13 we show a number of sequential recurrences for the population $\left|\psi_{0}(t)\right|^{2}$, where one clearly sees that the spikes gradually change from the bell-like (Gaussian) shape to a shape with asymmetric oscillations typical of the Airy function, since the relative role of the cubic term increases with time. Note that the spike duration and the period of spike oscillations changes in course of time as is clear from (3.95) where a slow decrease of the spike oscillation frequency with time is associated with the factor $\sqrt[3]{3 \alpha t}$ in the denominator of the argument of the Airy function. We also note that this universal Airy-like shape of the spikes is attained in all cases, no matter what the initial profile of the probability amplitude injection, provided the injection was short compared to the wavepacket oscillation period $T=2 \pi / \delta$ in the potential pit. Destruction of this picture occurs when the spike duration becomes comparable with the period $T$. In the wavepacket representation this means that at a certain moment the wavepacket becomes so broad that it spans the potential well. Corresponding temporal spikes start to overlap thus entering in the regime of strong interference, where the Poisson summation formula remaining formally valid does not yield a simple expression any longer. The separation of the probability amplitude into a sequence of


Fig. 3.14. (a) A series of population recurrences at longer time. The width of the spikes becomes comparable with the return period and the signal takes on a complicated shape. Corresponding wavepackets (b) span all of the potential pit.
individual signals no longer makes sense, and the corresponding population acquires a complex irregular shape resembling a random function as shown in Fig. 3.14.

### 3.3.3 Quantum Revivals

Now we show how the interference results in the complete reconstruction of the initial wavepacket, which is followed by the repetition of all the series
of recurrences. This phenomenon is called quantum revival. We dwell here mainly on the case of the regular, almost equidistant (that is slightly anharmonic) spectrum, which we were considering in the previous subsection, and discuss in just a few details another important case of an irregular spectrum, which will be needed later on in this chapter for qualitative estimates.

## Regular Spectrum

In order to get an insight into the origin of the quantum revival effect let us consider (3.86) in more detail. The cubic term in the exponent reads

$$
\begin{equation*}
i \alpha(n+1 / 2)^{3} t=i \alpha t\left(n^{3}+\frac{3}{2} n^{2}+\frac{3}{4} n+\frac{1}{8}\right) . \tag{3.97}
\end{equation*}
$$

For large $t$ this term is large and differs considerably for the neighboring $n$.


Fig. 3.15. Terms of the sum (3.86) as vectors in the complex plane. The corresponding Cornu constructions for the sum are shown in the inserts in the lower left corners. At $t=0$ all vectors are aligned and the modulus of the sum is maximum. With the course of time, each of the terms acquires a phase factor, such that the sum decreases. At a time equal to the return period, all of the vectors are almost aligned, since the phase difference comes only from the cubic anharmonicity. Complete alignment occurs at the time of complete revival, when the phase differences of the neighboring terms are multiples of $2 \pi$.

However, for some particular values of $\alpha t$ this big change becomes a multiple of $2 \pi i$ and therefore it does not result in any phase shift. Indeed, for $\alpha t=8 \pi$ it amounts to

$$
\begin{equation*}
i \alpha(n+1 / 2)^{3} t-i \alpha(n-1 / 2)^{3} t=i \pi\left(24 n^{2}+2\right) \tag{3.98}
\end{equation*}
$$

and yields zero phase shift for the amplitudes $\psi_{n}(t=8 \pi / \alpha)$ whatever the indices $n$ are. Note that in the case $\alpha /(\tau \delta)^{3} \ll \delta^{\prime}$ the cubic term (3.97)
changes negligibly if $t$ changes within one period $T$. This change can be ignored, the overall contribution of the anharmonicity $\alpha$ therefore vanishes, and we return to (3.87) for the probability amplitude of the state $|0\rangle$. This means that the sequence of Gaussian spikes (3.93), which has been destroyed by the anharmonicity $\alpha$ at the time $t>1 / \alpha(\tau \delta)^{3}$, appears again in its original form at the time $t \simeq 8 \pi / \alpha$. The whole scenario of the time evolution repeats again after this revival time $T_{r}=8 \pi / \alpha$. The wavepacket $\Psi\left(x, T_{r}\right)$ reconstructs itself back from the shape spread over all the potential pit (Fig. 3.14) to the initial Gaussian localized form (Fig. 3.12).

It is expedient to illustrate the process of recurrence and revivals with the help of vector diagrams in the complex plane, similar to the Cornu geometric construction (Cornu spiral) widely known in diffraction theory. Indeed, each term of the sum (3.86) is a complex number of modulus $e^{-(n+1 / 2)^{2} \delta^{\prime 2} \tau^{2} / 4}$ and phase $-i(n+1 / 2) \delta^{\prime} t+i \alpha(n+1 / 2)^{3} t$. We depict these numbers as vectors in the complex plane shown in Fig.3.15. All the vectors have zero phase at $t=0$ and therefore their sum has maximum absolute value $\psi_{0}(0)=$ 1. In the course of time the phase difference among the neighboring terms reduces the size of the absolute value of the sum, which almost vanishes at $t \sim \pi / \delta^{\prime}$, when the neighboring terms are opposite in phase. For the time scale where $\alpha t n^{3}$ is still a small value for typical $n \sim \delta^{\prime} \tau$, the sum again has a maximum at the multiples $t=T k$ of the recurrence time $T=2 \pi / \delta^{\prime}$, although this maximum is a bit smaller than the value at $t=0$ due to the influence of the anharmonicity $\alpha$, which introduce small phase shifts among the neighboring terms resulting in disalignment of the vectors representing these terms. The longer the time the stronger the disalignment, such that at $t \sim \delta^{\prime 3} / \alpha$ no revivals can be seen any longer. However, at the time $t \sim$ $T_{r v}=8 \pi / \alpha$ the phase differences resulting from the anharmonicity become multiples of $2 \pi$, and therefore complete revivals become again possible. In


Fig. 3.16. Evolution at the initial stage (a) and quantum revival (b) at $t=T_{r v}=$ $10^{4} T$. The number of populated levels is $2 / \delta^{\prime} \tau=40$. The vectors representing the terms of the sum ( 3.86 ) become aligned again, as at the beginning, and the evolution repeats.

Fig. 3.16 we show the probability $\rho$ to be in the state $|0\rangle$ at the initial stage of the evolution at $t \sim 0$ and at the time $t \sim T_{r v}$ of the revival. One sees the complete reconstruction of the wavepacket at $t=T_{r v}$ followed by the repetition of the evolution. Note that here we have considered the case of a cubic nonlinearity, although the same phenomenon of revivals takes place also for a square nonlinearity $\alpha t n^{2}$ for a different $T_{r v}$.

## Irregular Spectrum

The important case of the irregular spectrum of eigenstates should also be mentioned here. It is rather typical for a multidimensional motion and does not result in well-pronounced recurrences after a time of the order of the inverse density of states. Indeed, in classical mechanics the motions in different directions are usually nolredyseparable for two or more dimensional finite systems of general type, as we have already discussed for Sinai billiards (Fig. 1.1). This implies that the classical particle never returns exactly to the initial position, although, according to the Pointcaré theorem, it can approach its vicinity with any predetermined accuracy. The smaller the predetermined distance the longer the time needed to get in the vicinity. On the contrary, a finite size of wavepackets in quantum mechanics implies that during such a return a part of the population comes back to the initial state, which means a partial recurrence. Since a long classical time is typically required for return at the vicinity of a size of the initial wavepacket, the packet dispersion during this time is large, and therefore the recurrency amplitude is typically small, resembling rather the regime shown for one-dimensional motion in Fig. 3.14.

In the same time an analog of complete or partial revival can occur in such a system. Revivals of amplitudes of the order of unity are extremely rare events, which resemble giant fluctuations. Let us estimate the time needed to return back to the initial state with an accuracy $\epsilon$. By analogy to (3.86) we take the amplitude in the form

$$
\begin{equation*}
\psi_{0}(t) \sim \sum_{n} e^{-(n+1 / 2)^{2} / N^{2}-i \Delta_{n} t} \tag{3.99}
\end{equation*}
$$

where $N$ stands for the typical number of the populated energy eigenstates and $\Delta_{n}$ denote detunings of these states from the energy of the level. Actually a simpler expression

$$
\begin{equation*}
\psi_{0}(t)=\frac{1}{2 N+1} \sum_{n=-N}^{N} e^{-i \Delta_{n} t} \tag{3.100}
\end{equation*}
$$

is sufficient for the estimates, where we have replaced the bell-shaped Gaussian distribution by a rectangular distribution, which affects only the shape of revival signals but not the amplitudes at maximum. In order to have the sum (3.100) close to unity we should find a time $t$ such that all the phases $\Delta_{n} t$ are
close to multiples of $2 \pi$. This means that the end of the $(2 N+1)$-dimensional vector

$$
\vec{r}(t)=\left(\Delta_{-N} t, \Delta_{-N+1} t, \ldots ; \Delta_{n} t, \ldots \Delta_{N} t\right)
$$

should be inside of a sphere of radius $\sqrt{2 \epsilon N}$ around the point

$$
P=2 \pi\left(m_{-N}, m_{-N+1}, \ldots ; m_{n}, \ldots m_{N}\right)
$$

with integer $m_{n}$. If all $\Delta_{n}$ are random, one can identify the vector $r(t)$ with coordinates of a particle moving with a random velocity

$$
\vec{v}=\left(\Delta_{-N}, \Delta_{-N+1}, \ldots ; \Delta_{n}, \ldots \Delta_{N}\right)
$$

in the $(2 N+1)$-dimensional space and the $\epsilon$ vicinities of points $P$ with solid spherical scatterers of radii $\sqrt{2 \epsilon / N}$ placed in the nodes of a $2 \pi$-periodic lattice. By analogy to molecular physics it allows one to say that the revival mean free time $T_{r v}$ equals the inverse relative cross-section $\Gamma(N) \epsilon^{-N} \pi^{-N}$ of the $(2 N+1)$-dimensional spherical scatterer. With the help of the Stirling formula for the $\Gamma$-function we find

$$
\begin{equation*}
T_{r v} \sim e^{-N}(N / \epsilon \pi)^{N} \tag{3.101}
\end{equation*}
$$

which shows that the revival time rapidly increases with the number $N$ of the energy eigenstates involved. Indeed, the smaller the total mismatch $\epsilon$ of the revival is, the smaller should be the mismatches $\epsilon / N$ of the individual terms of the sum (3.100). This corresponds to the phase mismatches $\sim \sqrt{\epsilon / N}$ at the position, when the vectors representing the terms are almost aligned and $\cos (x) \simeq 1-x^{2} / 2$. Probability to have all $2 N$ terms aligned with this accuracy is of the order of $(\sqrt{\epsilon / N})^{2 N}$, which yields (3.101) apart of a numerical factor $1 / \sqrt{e \pi}$ for each term of the sum.

### 3.3.4 Fractional Revivals

Let us now return to the almost equidistant spectra. One observes regular almost periodic time dependence of the probability amplitude $\psi_{0}(t)$ not only at the initial stage of the evolution and at the time $T_{r v}=8 \pi / \alpha$ of the complete revival, but also at times $t \sim T_{f r}=(p / q) T_{r v}$, where the fraction is given by the ratio mutually prime numbers $p$ and $q$. This phenomenon is known as Averbukh-Perelman fractional revivals, and it results in a more complicated shapes of $\psi_{0}(t)$ as compared to the simple revivals. In order to be consistent with the consideration of previous sections we concentrate here on the case of cubic anharmonicity, although the fractional revivals phenomenon can also be observed for quadratic nonlinearity. The wavepacket picture corresponding to the fractional revivals shows that the particle usually is distributed among two or several well-localized wavepackets each of which contains just a fraction of the particle. In other words, the wavefunction $\Psi\left(x, T_{r}\right)$ is a quantum


Fig. 3.17. Fractional revivals at $t=T_{r v} / 12=40000 T / 12$. This regime corresponds to complete revival that, however, occurs not at the multiple of the return period. The number of the populated levels is $2 / \delta^{\prime} \tau=18$. (a) The time dependent population of the initial state. (b) The wavepacket representation: $\rho=|\Psi(x, t)|^{2}$ at $\tau=3333.333$. Only one wavepacket is seen in the potential pit. (c) The ends of vectors representing the terms of the sum (3.86) are shown schematically by points in the complex plane for corresponding times. The revival corresponds to the complete alignment of all the terms.
superposition of several well-separated Erenfest wavepackets, which allows us to consider this quantum mixture of several classical states as a Schrödinger cat. For some $p$ and $q$ the wavefunction may consist of a single wave packet which returns to the initial state not at a multiple of the return time $n T$ as it was at the beginning, but at a moment $n T+\Theta T$, where $\Theta T$ is a fraction $\Theta$ of the return period $T$. We depict such a situation for $p=1 ; q=12$ in Fig. 3.17.

c)


Fig. 3.18. Same as Fig. 3.17 for a more complex situation at $t=T_{r v} / 24$. The number of the populated levels is $2 / \delta^{\prime} \tau=18$ as before. (a) The time-dependent population of the initial state has two spikes. (b) Two counterpropagating wave packets are seen at $\tau=1666.56$ in the potential pit. (c) These fractional revivals correspond to the situation where the terms of the sum (3.86) are grouped in two sets of aligned vectors, with a $\pi / 2$ phase shift between the sets.

For a shorter time ( $p=1 ; q=24$ ) one observes two wavepackets and two corresponding spikes of height $1 / 2$. We depict such a situation in Fig. 3.18. At the time points corresponding to the position of fractional revivals terms of the sum (3.86) are grouped into two sets with a mutual phase shift of $\pi / 2$, as one can see from panel (c).

A more complicated situation takes place at $T_{f r} \sim T_{r v} / 32$, when four spikes of different size are seen during one return period. We depict such a situation in Fig. 3.19.


Fig. 3.19. Same as Fig. 3.17 for $t=T_{r v} / 32$. (a) The time-dependent population of the initial state has four spikes of different size. (b) Four counterpropagating wavepackets of different height corresponding to these spikes are seen at $\tau=1250.333$. (c) Grouping of the terms of (3.86) to aligned sets is less apparent.

The grouping of vectors representing the terms of (3.86) to aligned sets is less evident in this case. However they yield four counterpropageting Gaussian wave packets for the equivalent particle in the potential pit. One can conclude that in the regime $\alpha T\left(1 / \delta^{\prime} \tau\right)^{3} \ll 1$ under consideration the spikes of the fractional revivals have the same shape as the initial signal, but the probability of return is shared among two or several time intervals relatively short as compered to the return period $T$.

Let us dwell now on the distribution law for the return probability amplitudes among the spikes. To this end we consider terms of the sum (3.86) at $t=(p / q) T_{r v}+\vartheta=(8 \pi p / \alpha q)+\vartheta$, where $0<\vartheta<T$ is a time interval of the
duration of one return period. The general term of the sum adopts the form

$$
\begin{align*}
& \exp \left[-(n+1 / 2)^{2} \frac{\delta^{\prime 2} \tau^{2}}{4}+i\left(\frac{p \pi}{q}+\frac{6 n p \pi}{q}+\frac{12 n^{2} p \pi}{q}+\frac{8 n^{3} p \pi}{q}-\frac{4 p \pi \delta^{\prime}}{q \alpha}\right.\right. \\
& \left.\left.-\frac{8 n p \pi \delta^{\prime}}{q \alpha}+\frac{\alpha \vartheta}{8}+\frac{3 n \alpha \vartheta}{4}+\frac{3 n^{2} \alpha \vartheta}{2}+n^{3} \alpha \vartheta-\frac{\delta^{\prime} \vartheta}{2}-n \delta^{\prime} \vartheta\right)\right] . \tag{3.102}
\end{align*}
$$

We neglect the terms proportional to $\alpha \vartheta$, since they are small in the regime under consideration, omit the multiples of $2 \pi$, and arrive at

$$
\begin{equation*}
\exp \left\{-(n+1 / 2)^{2} \frac{\delta^{\prime 2} \tau^{2}}{4}+i\left[\frac{12 n^{2} p \pi}{q}+\frac{8 n^{3} p \pi}{q}+n\left(\frac{12 p \pi}{q}-\frac{8 p \pi \delta^{\prime}}{q \alpha}-\delta^{\prime} \vartheta\right)\right]\right\} \tag{3.103}
\end{equation*}
$$

Let us denote by $\Theta$ the fractional part of $\left(\frac{12 p \pi}{q}-\frac{8 p \pi \delta^{\prime}}{q \alpha}-\delta^{\prime} \vartheta\right) / 2 \pi$, which allows one to express a moment of time in the interval between two sequential return periods $l T$ and $(l+1) T$ in terms of a fraction of $T$. Then we set $q=2 q^{\prime}$ and obtain the sum (3.86) in the form

$$
\begin{equation*}
\sum_{n} \exp \left\{-(n+1 / 2)^{2} \frac{\delta^{\prime 2} \tau^{2}}{4}+i\left(\frac{6 n^{2} p \pi}{q^{\prime}}+\frac{4 n^{3} p \pi}{q^{\prime}}+n 2 \pi \Theta\right)\right\} \tag{3.104}
\end{equation*}
$$

We note that the parts of the phase which depend on $n^{3}$ and $n^{2}$ yield multiples of $2 \pi$ for $n=m q$, where $m$ is an integer. We can therefore set $n=m q^{\prime}+k$ and replace the summation over $n$ by two sums:

$$
\begin{align*}
& \sum_{m=-\infty}^{m=\infty} \sum_{k=1}^{k=q^{\prime}} \exp \left\{-\left(m q^{\prime}+k+1 / 2\right)^{2} \delta^{\prime 2} \tau^{2} / 4\right. \\
& \left.+i\left[\frac{6\left(m q^{\prime}+k\right)^{2} p \pi}{q^{\prime}}+\frac{4\left(m q^{\prime}+k\right)^{3} p \pi}{q^{\prime}}+2 \pi\left(m q^{\prime}+k\right) \Theta\right]\right\} \tag{3.105}
\end{align*}
$$

We replace summation over $m$ by integration over $x$ with the help of the Poisson formula $\sum_{m} \delta(x-m)=\sum_{s} e^{i 2 \pi s x}$ (3.89) and perform this Gaussian integration, which yields

$$
\begin{equation*}
\sum_{s=-\infty}^{\infty} \sum_{k=1}^{q^{\prime}} \frac{2 \sqrt{\pi}}{q^{\prime} \delta \tau} e^{i \pi\left(6 k^{2} p / q^{\prime}+4 k^{3} p / q^{\prime}-2 k s / q^{\prime}\right)-i \pi\left(s / q^{\prime}+\Theta\right)-\frac{4 \pi^{2}}{\delta^{2} \tau^{2}}\left(s / q^{\prime}+\Theta\right)^{2}} \tag{3.106}
\end{equation*}
$$

It is convenient to rewrite this equation in the form

$$
\begin{equation*}
\sum_{s=-\infty}^{\infty} \frac{2 \sqrt{\pi}}{q \delta \tau} e^{-i \pi\left(s / q^{\prime}+\Theta\right)-\frac{4 \pi^{2}}{\delta^{2} \tau^{2}}\left(s / q^{\prime}+\Theta\right)^{2}} \sum_{k=1}^{q^{\prime}} e^{i \pi\left(6 k^{2} p / q^{\prime}+4 k^{3} p / q^{\prime}-2 k s / q^{\prime}\right)} \tag{3.107}
\end{equation*}
$$

where one clearly identifies normalized Gaussian spikes with zero phase at the maximum given by terms of the first sum at times $T \Theta$ with $\Theta=-s / q^{\prime}$ and a universal amplitude function

$$
\begin{equation*}
\operatorname{Fr}\left(p, q^{\prime}, s\right)=\sum_{k=1}^{q^{\prime}} e^{i \pi\left(6 k^{2} p / q^{\prime}+4 k^{3} p / q^{\prime}-2 k s / q^{\prime}\right)} \tag{3.108}
\end{equation*}
$$



Fig. 3.20. A universal dependence $\operatorname{Fr}\left(p, q^{\prime}, s\right)$ which determines the distribution of the amplitudes of the fractional revivals of population as a function of the fractional revival number $s$ for the case of cubic anharmonicity $\alpha(n+1 / 2)^{3}$. Normally the fractional revivals in the time region $t \sim \frac{p}{2 q^{\prime}} \frac{8 \pi}{\alpha}$ occur at the moments timeof $t=$ $l T+\frac{s}{q^{\prime}} T$, although for some $s$ the revival amplitude may vanish, depending on the factorization properties of the numbers $p$ and $q=2 q^{\prime}$.
depending only on the integer numbers, enumerating the revival time $(s)$ and the revival type $(p, q)$.

For each type of nonlinear dependence of the eigenstate energy on the number $n$, the function $\operatorname{Fr}\left(p, q^{\prime}, s\right)$ of integer variables $p, q^{\prime}$, and $s$ is determined by (3.108), and yields the probability of fractional revivals at the time $t=l T+s t / q^{\prime}$ for the time region $t \sim l T \sim p T_{r v} / 2 q^{\prime}$. Note, that an integer $l$ enumerating the number of the return period for which the fractional revivals take place apparently equals the integer part of the ratio $p T_{r v} / 2 q^{\prime} T$. For the case of cubic anharmonicity under consideration the function $\operatorname{Fr}\left(p, q^{\prime}, s\right)$ is depicted in Fig. 3.20 for several lowest values of $q$. One sees that for the prime $q^{\prime}=5,7,11$ all the fractional revivals at different $s$ are present although with different amplitudes, whereas for composed $q^{\prime}=4,6,8,9,10$ some of the spikes disappear. This indicates a hidden relation between the fractional revivals and number theory.

We have considered here the quantum revivals in the simplest version, when all wavepackets were initially of a bell-like Gaussian shape, and we have called revivals and fractional revivals the restoration of such a shape at a given time. This does not always mean that the quantum system returns completely to the initial level, but only implies that the initial state population time dependence $\rho_{0}(t)$ is restored at some moment as it was at the
beginning, that is in the course of the population injection. For the original level-band problem, with the initial condition $\psi_{0}(t=0)=1$, when the population does indeed return completely back to the initial state, the wavepackets have different forms - the initially exponential shape gradually transforms to an Airy function-like dependence. In such a situation the variety of possible shapes of fractional revivals is richer, although apart of this circumstance the phenomenon keeps the same general features given by the function $\operatorname{Fr}\left(p, q^{\prime}, s\right)$. We also note that for any other type of nonlinearity, like $\alpha n^{2}, \alpha n^{4}, \ldots$ etc., the corresponding function $\operatorname{Fr}\left(p, q^{\prime}, s\right)$ is different and can be constructed by analogy.

### 3.3.5 Revivals and the Classical Limit

The existence of revivals and fractional revivals suggests to consider more carefully the traditional classical limit of quantum mechanics, at least in the case of the one-dimensional motion where the energy eigenstate spectrum demonstrates a tendency to become equidistant. Indeed, the parameter $n=$ $S / \hbar$ responsible for the transition to the classical limit is the ratio of the mechanical action $S=\int_{\text {period }} p(q) d q$ of the particle and the elementary action given by the Planck constant $\hbar$. It gives the mean $\bar{n}$ of the typical number $n$ of the populated quantum states. The larger the ratio $S / \hbar$, the more classical the system is. The transition frequency $\omega_{n}=\left(E_{n+1}-E_{n}\right) / \hbar$ among the neighboring states $n$ and $n+1$ is a frequency of classical periodic motion, and it is apparently a classical quantity which may depend on mass and elastic constant but not on $\hbar$. At the same time, the quantum anharmonicity of the motion, that is the ratio of the anharmonic correction $\omega_{n+1}-\omega_{n}$ of the transition frequency for the adjacent pairs of the neighboring quantum levels to the frequency $\omega_{n}$ itself, scales as $\hbar / S$. Therefore with increasing $n$ the one-dimensional system becomes more and more harmonic.

The classical limit of quantum mechanics is based on the Erenfest theorem, which states that the center of a wavepacket composed as a linear combination of many energy eigenstates follows the classical trajectory. For the harmonic oscillator this is always the case. However the majority of the one-dimensional systems are anharmonic, which results, in particular, in a spread of the wavepacket increasing in the course of time. This quantum dispersion destroys the localization of the Erenfest wavepackets. Hence the more classical the system the more equidistant the spectrum is, and the longer is the relative time needed for the spread of a wavepacket composing a localized particle.

In the traditional classical limit of quantum mechanics one first takes $\hbar=0$ and then considers the limit $t \rightarrow \infty$, which implies the time domain of interest remaining shorter than the wavepacket dispersion time $\sim S / \hbar \omega_{n}$, while implementing the consistent quantum approach one notes that at the long time scale $t \gg S / \hbar \omega_{n}$ the anharmonic terms are important, since they wash out the sharp distribution of the quasiclassical particle on the classical
trajectory in the manner that the particle becomes distributed along all the length of the trajectory with almost uniform probability, which depends only on the inverse residual time at a chosen point. Figure 3.14, which corresponds to the same situation for a relatively small number of energy eigenstates in the wavepacket, gives an idea about the situation in the classical limit.

From the example of revivals and fractional revivals we have learned that the consistent quantum consideration of the particle propagation problem yields another regime, where the particle is localized in one or a few wavepackets obeying classical mechanics, that is the particle spread all over the trajectory in the course of evolution reinstates itself again as a single welllocalized wavepacket or as a quantum superposition of the well-separated and well-localized wave packets, each of which follows the same classical trajectory just being displaced in time by a fraction of the oscillation period. In the classical limit such a regime corresponds to an intermediate asymptotic $\hbar \rightarrow 0, t \rightarrow \infty$, but with the constraint $\hbar t=$ const for the case of quadratic anharmonicity $\alpha n^{2}$. The size of the constant is of the order of a typical action $S$ multiplied by the oscillation period $T=2 \pi / \omega_{n}$. For cubic anharmonicity $\alpha n^{3}$ the ratio $\left(\omega_{n+1}-\omega_{n}\right) / \omega_{n}$ scales as $(\hbar / S)^{2}$ and this constraint reads $\hbar^{2} t=$ const, with the constant of the order of $S^{2} T$.

The considered asymptotic implies that during the time of evolution the quantum system remains isolated, that is it obeys Hamiltonian mechanics. For really classical objects the condition $\hbar t=S T$ is difficult to fulfill, since the typical actions are large, and the typical periods are long. Isolating a system for such a long period of time as $t=S T / \hbar$ is unrealistic. But for the case of mesoscopic systems, such as cold atoms or molecules at relatively large distances, this condition can in principle be satisfied.

### 3.4 Population of Inhomogeneous Bands

The case of an equidistant spectrum coupled to a level by identical matrix elements is one of a few examples of the level-band problem that can be solved exactly and analytically in the traditional manner. It allows to get explicit analytical expressions for a number of different interference effects that may occur in the system and that have been discussed earlier in this chapter. However, no exact analytical solution is available for a more complex quantum system, when the coupling constant $V_{0 n}$ differs significantly with the band level number $n$, and the detunings $\Delta_{n}$ of the band levels deviate from the equidistant positions $\delta n$. Of course, the exact expression (3.26) can still be written, but neither the inverse Fourier transformation can be carried out exactly, nor an adequate analysis of this exact expression is possible. We encounter the situation, discussed earlier in Chap. 2, where the idea of the ensemble average becomes indispensable. It yields an analytical expression of the population of the level, although not for a particular quantum levelband system, but for the population averaged over many possible realizations
of such a system. We note that in such systems one can encounter nonexponential behavior which is closely related with essentially non-Wigner level statistics.

### 3.4.1 Statistically Independent Levels

Let us assume that we have an ensemble of level-band systems, each of which has exactly the same numbers $N\left(V_{i}\right)$ of band levels coupled to the state $|0\rangle$ by the matrix elements $V_{0 n}=V_{i}$. Although the numbers of levels with a given coupling $V_{0 n}$ are taken to be the same for each system comprising the ensemble, the values of the couplings $V_{i}$ may change significantly, that is by orders of magnitude. The positions of the levels, on the contrary, are different for different systems, such that the detuning $\Delta_{n}$ may take any value from a broad strip $\Delta_{n} \in(-\Gamma, \Gamma)$. An independent distribution of the level energies is called "Poissonian statistics". We assume that the width of the strip $\Gamma \rightarrow \infty$, such that the mean density $g_{i}=N\left(V_{i}\right) / 2 \Gamma$ of the levels with a fixed coupling remains constant. An example of such a system is shown in Fig. 2.5.

How many band levels are in resonance with the state $|0\rangle$ in such a situation? From the discussion of Sect.3.2.2 we remember that a two-level system is resonant when the detuning between the levels is less than their coupling. Let us calculate the probability $P_{r}$ to have a resonance suggested by this criterion. Consider all the levels $n$ of type $i$ that have the matrix elements of coupling $V_{0 n}=V_{i}$. The ratio of the coupling and a typical smallest detuning $\Delta^{\prime} \sim 1 / g_{i}$ of the level of the group from resonance can be estimated by the product $g_{i} V_{i}$. As long as $g_{i} V_{i} \ll 1$, this ratio gives the probability to have a resonance in the selected group of levels. In order to get the total probability to have a resonance one has to sum up all the types of level and arrive at $P_{r}=\sum_{i} g_{i} V_{i}$. In a system with distributed parameters the sum has to be replaced by the integral

$$
\begin{equation*}
P_{r}=\int V g(V) d V \tag{3.109}
\end{equation*}
$$

where $g(V) d V$ is a density of states with $V_{0 n} \in(V, V+d V)$.
While $P_{i}$ is small it does indeed yield the probability to have a resonance. However, the integral (3.109) may be large or may even diverge at the lower limit, which means that the expected number of resonances is much larger than one or is even infinite. Then the probability $P_{i}$ suggested by the criterion for a two-level system does not represent the expectation value of the number of resonances, since in a multilevel system the resonance criterion is different. From the discussion of Sect. 3.2.2 we know that $N_{R} \sim 2 \pi P_{r}^{2}$ gives a better estimate of the number $N_{R}$ of resonant levels, provided the integral (3.109) converges.

If the integral (3.109) diverges at the lower limit, the largest contribution is given by the weakly coupled levels. This results in a very special dynamics of the level-band interaction which implies that the number of levels involved


Fig. 3.21. The probability amplitude (a) and the probability (b) to be in the state $|0\rangle$ for a level-band system with coupling changing by orders of magnitude, shown in the logarithmic time-scale. Complete revivals become more and more seldom with the course of time, and therefore the time average population decreases. The ensemble average washes out the oscillations and yield zero average amplitude while the average population decreases as a power of the time, either $\rho_{0}(t) \sim t^{2(\alpha-2)}$ for $\alpha<2$ or $\rho_{0}(t) \sim t^{-2(\alpha-2) /(\alpha-1)}$ for $\alpha>2$, where $g(V) \sim g_{\alpha} V \alpha$ is the asymptotic behavior of the state density as $V \rightarrow 0$.
in the process increases in the course of time. At the first stage of the process the population rapidly oscillates among a relatively small number of states for which $V t \sim 1$. In the course of time, more and more weakly coupled levels satisfy this condition and therefore become important, and hence a number of slow Rabi oscillations replace the rapid process. This means that the typical frequency scales down in the course of time, until it becomes zero. One gets an idea about such a regime by considering a function $\phi(a, t)=\prod_{n=0}^{\infty} \cos \left[t e^{-a n}\right]$, where harmonic oscillations at unit frequency $\cos [t]$ are modulated by a series of periodically changing amplitudes $\cos \left[t e^{-a n}\right]$ of exponentially decreasing frequencies $e^{-a n}$. In Fig. 3.21 we show the probability amplitude and the probability suggested by the model $\psi_{0}(t)=\phi(a, t)$. The oscillation frequencies are different for different systems comprising the ensemble, and therefore the ensemble average probability amplitude is not a self-averaging value and already vanishes for relatively short evolution times. But the ensemble average population is well defined, and therefore considering the essentially non-homogeneous bands we concentrate on this very quantity.

Before discussing the technique of the ensemble average it is expedient to consider another aspect of the problem, and turn to the location of the energy eigenvalues by analogy to the approach shown in Fig. 3.11. The divergency of the integral (3.109) at the lower limit means, that the poles of the sum $\sum_{n} V_{0 n} V_{n 0} /\left(\varepsilon-\Delta_{n}+i 0\right)$ in (3.26), which for the equidistant spectrum were placed regularly in the real axis thus giving rise to the cotangent function, are now randomly spread along this axis in such a manner that an infinite number of singularities is present around each real point. Clearly the same is also the case for the zeros of the denominator of (3.26), or, in other words,
each point of the real axis is an essential singularity of the Fourier transform $\psi_{0}(\varepsilon)$ of the type known as a limiting point of poles. The density of poles increases to infinity while their strength $V_{0 n} V_{n 0}$ decreases. The multitude of the pole positions resembles in some respect the Cantor set.

### 3.4.2 Factorization of the Level Population and the Ensemble Average

We show now how to perform the ensemble average. The Schrödinger equation for a level-band system in the Fourier representation yields (3.26) in the form

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{i}{\varepsilon-\sum_{n} V_{0 n} V_{n 0} /\left(\varepsilon-\Delta_{n}+i 0\right)} \tag{3.110}
\end{equation*}
$$

In the situation where all the levels are similar and each of them is weakly populated and hence never absorbs a significant fraction of the total population, one can replace the sum by the integral $\sum_{n} \ldots \rightarrow \int g(\Delta, V) d \Delta d V \ldots$ at $t \ll g$ as we have done considering the continuous-band model. Here $g(\Delta, V) d \Delta d V$ denotes the expectation number of levels in an interval $d \Delta$ around $\Delta$ that is coupled to the state $|0\rangle$ by the matrix element in the interval $d V$ around $V$. This replacement is illegal for matrix elements $V_{0 n}$ that are different by orders of magnitude, since different groups of band levels (sorted by the size of $V_{0 n} V_{n 0}$ ) result in recurrences, revivals and fractional revivals in essentially different time-scales, which also means that some of the levels can accumulate a significant fraction of the total population more often then the other, and hence during a finite interval of time they contribute to the sum (3.110) with higher weights. In such a situation the statistical spread of various physical values is large, and in order to get a consistent description one has to perform the ensemble average of the quantity of interest, that is for the population. To this end we take the exact expression for the population

$$
\begin{align*}
\rho_{0}(t)= & \psi_{0}(t) \psi_{0}^{*}(t) \\
= & \frac{1}{4 \pi^{2}} \int \frac{e^{-i(\varepsilon-\xi) t}}{\varepsilon-\sum_{m} V_{1 m} V_{m 1} /\left(\varepsilon-\Delta_{m}+i 0\right)} \\
& \frac{1}{\xi-\sum_{m} V_{1 m} V_{m 1} /\left(\xi-\Delta_{m}-i 0\right)} d \varepsilon d \xi \tag{3.111}
\end{align*}
$$

and perform the ensemble average over $\prod d \Delta_{n} d V_{0 n}$ with a distribution function $g\left(\left\{\Delta_{n}\right\},\left\{V_{0 n}\right\}\right)$ depending the sets of all detunings $\left\{\Delta_{n}\right\}$ and all couplings $\left\{V_{0 n}\right\}$.

Our assumption implies independent (Poisson) statistics of the position of the band states and their couplings to the level, that is $g\left(\left\{\Delta_{n}\right\},\left\{V_{0 n}\right\}\right)=$ $\prod g\left(\Delta_{n}\right) g\left(V_{0 n}\right)$, with $g\left(\Delta_{n}\right)=\Theta\left(\Gamma-\left|\Delta_{n}\right|\right) / 2 \Gamma$ given by the step function $\Theta(x)$. However, even within so strong a simplifying approximation the
ensemble average procedure still requires some effort. For this purpose, one first has to represent the right-hand side of (3.111) in factorized form, that is as a product of terms, each of which depends only on the detuning $\Delta_{n}$ and coupling $V_{0 n}$ of the same level $n$. Then the ensemble average is reduced to an independent average over the position and the coupling of each level.

Let us first take the average of the population (3.111) over the detunings. It reads

$$
\begin{align*}
\left\langle\rho_{0}(t)\right\rangle= & \frac{1}{4 \pi^{2}} \prod_{n}\left\{\int_{-\Gamma}^{\Gamma} \frac{d \Delta_{n}}{2 \Gamma}\right\} \int_{-\infty}^{\infty} d \varepsilon d \xi \\
& \times \frac{e^{-i(\varepsilon-\xi) t}}{\left(\varepsilon-\sum_{n} \frac{V_{0 n} V_{n 0}}{\left(\varepsilon-\Delta_{n}+i 0\right)}\right)\left(\xi-\sum_{n} \frac{V_{0 n} V_{n 0}}{\left(\xi-\Delta_{n}-i 0\right)}\right)} \tag{3.112}
\end{align*}
$$

We introduce new variables $\zeta=\varepsilon-\xi, \eta=(\varepsilon+\xi) / 2, x_{n}=\Delta_{n}-\eta$ and arrive at

$$
\begin{align*}
\left\langle\rho_{0}(t)\right\rangle= & \frac{1}{4 \pi^{2}} \prod_{n}\left\{\int_{-\Gamma-\eta}^{\Gamma-\eta} \frac{d x_{n}}{2 \Gamma}\right\} \int_{-\infty+i \nu}^{\infty+i \nu} \int_{-\infty}^{\infty} d \zeta d \eta e^{-i \zeta t}  \tag{3.113}\\
& \times \frac{1}{\left(\eta+\zeta / 2-\sum_{n} \frac{\left|V_{0 n}\right|^{2}}{\left(\zeta / 2-x_{n}+i 0\right)}\right)\left(\eta-\zeta / 2-\sum_{n} \frac{\left|V_{0 n}\right|^{2}}{\left(-\zeta / 2-x_{n}-i 0\right)}\right)},
\end{align*}
$$

where the integration contour for $\zeta$ is the same as for $\varepsilon$ (see $C_{1}$ in Fig. 3.2), and the integral over $d \eta$ is calculated along the real axis. We now neglect $\eta$ with respect to large $\Gamma$ in the limits of integration over $d x_{n}$ and perform integration over $d \eta$, which yields

$$
\begin{equation*}
\left\langle\rho_{0}(t)\right\rangle=\int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \zeta}{2 \pi i} \prod_{n}\left\{\int_{-\Gamma}^{\Gamma} \frac{d x_{n}}{2 \Gamma}\right\} \frac{e^{-i \zeta t}}{\zeta-\sum_{n}\left|V_{0 n}\right|^{2}\left(\frac{1}{\zeta / 2-x_{n}}+\frac{1}{\zeta / 2+x_{n}}\right)} \tag{3.114}
\end{equation*}
$$

Note that on the integration path in the $\zeta$-plane the denominator of the integrand (3.114) has a positive imaginary part, and hence the imaginary part of the fraction is negative. We can therefore make use of the representation $1 / x=i \int_{0}^{\infty} e^{i x \tau} d \tau$ and write

$$
\begin{align*}
\frac{1}{\zeta-\sum_{n} \frac{\left|V_{0 n}\right|^{2} \zeta}{(\zeta / 2)^{2}-x_{n}^{2}}} & =i \int_{0}^{\infty} \exp \left[i \tau \zeta-\sum_{n} \frac{\left|V_{0 n}\right|^{2} \zeta \tau i}{(\zeta / 2)^{2}-x_{n}^{2}}\right] d \tau \\
& =i \int_{0}^{\infty} e^{i \tau \zeta} \prod_{n}\left\{\exp \left[\frac{-i\left|V_{0 n}\right|^{2} \zeta \tau}{(\zeta / 2)^{2}-x_{n}^{2}}\right]\right\} d \tau \tag{3.115}
\end{align*}
$$

This representation gives the intended factorized form of the expression and allows one to separate the variables and perform the integration over $d x_{n}$ independently for each $n$. Substitution of (3.115) into (3.114) yields

$$
\begin{equation*}
\left\langle\rho_{1}(t)\right\rangle=\int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \zeta}{2 \pi} e^{-i \zeta t} \int_{0}^{\infty} d \tau e^{i \tau \zeta} \prod_{n}\left\{\frac{1}{2 \Gamma} \int_{-\Gamma}^{\Gamma} d x_{n} \exp \left[\frac{-i\left|V_{0 n}\right|^{2} \zeta \tau}{(\zeta / 2)^{2}-x_{n}^{2}}\right]\right\} \tag{3.116}
\end{equation*}
$$

In the limit under consideration of a large number of states and a large width of the band $\Gamma$ (with a constant state density) the value $\kappa=$ $\frac{1}{2 \Gamma} \int_{-\Gamma}^{\Gamma} d x \exp \left[-4 i\left|V_{0 n}\right|^{2} \zeta \tau /\left(\zeta^{2}-4 x^{2}\right)\right]$ is very close to unity, since the expression in the exponent vanishes for large $x$, and hence the integrand tends to unity almost everywhere within the integration limits. We make use of this fact and replace $\kappa$ by $e^{\kappa-1}$, which for (3.116) yields

$$
\begin{align*}
\left\langle\rho_{1}(t)\right\rangle= & \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \zeta}{2 \pi} e^{-i \zeta t} \int_{0}^{\infty} d \tau e^{i \tau \zeta} \\
& \exp \left\{\sum_{n}\left[\frac{1}{2 \Gamma} \int_{-\Gamma}^{\Gamma} \exp \left(\frac{4 i\left|V_{0 n}\right|^{2} \zeta \tau}{4 x^{2}-\zeta^{2}}\right) d x-1\right]\right\} \tag{3.117}
\end{align*}
$$

We note that $\exp \left[4 i\left|V_{0 n}\right|^{2} \zeta \tau /\left(4 x^{2}-\zeta^{2}\right)\right]$ is always 1 along the contour $C_{2}$ shown in Fig. 3.2, and hence the expression in the square brackets in (3.117) can be considered as a single integral along the contour $C=C_{1}+C_{2}$ of Fig. 3.2, with the part $C_{1}$ coinciding with the real axis:

$$
\begin{equation*}
\int_{-\Gamma}^{\Gamma} d x\left[\exp \left\{\frac{4 i\left|V_{0 n}\right|^{2} \zeta \tau}{4 x^{2}-\zeta^{2}}\right\}-1\right]=\int_{C} d x \exp \left\{\frac{4 i\left|V_{0 n}\right|^{2} \zeta \tau}{4 x^{2}-\zeta^{2}}\right\} \tag{3.118}
\end{equation*}
$$

In this limit one can also replace $\sum_{n} 1 / 2 \Gamma$ in (3.117) by the integral over $g(V) d V$. Note that singularities of the integrand are at different parts of the complex $x$-plane divided by the real axis, and hence the contour $C$ can be transformed to a small loop around one of these singularities - the point $x=\zeta / 2$ at the upper part. After replacement $x \rightarrow[1+2 /(x-1)] \zeta / 2$ the integral takes the form

$$
\begin{equation*}
I=-\int_{C} \frac{\zeta}{(x-1)^{2}}\left(\exp \left\{i \frac{\left|V_{0 n}\right|^{2} \tau}{\zeta} \frac{(x-1)^{2}}{x}\right\}\right) d x \tag{3.119}
\end{equation*}
$$

and the integration contour $C$ becomes a clockwise oriented circle of large radius centered at the point $x=1$. We now perform the integration iby parts around this point, which results in

$$
\begin{align*}
I & =\int_{C} \frac{\zeta}{1-x} \frac{d}{d x}\left(\exp \left\{i \frac{\left|V_{0 n}\right|^{2} \tau}{\zeta}\left(x-2+\frac{1}{x}\right)\right\}\right) d x \\
& =-i\left|V_{0 n}\right|^{2} \tau \int_{C} \frac{x+1}{x^{2}}\left(\exp \left\{i \frac{\left|V_{0 n}\right|^{2} \tau}{\zeta}\left(x-2+\frac{1}{x}\right)\right\}\right) d x \tag{3.120}
\end{align*}
$$

For this integral the radius of the circle $C$ can be taken equal to unity, and therefore by changing variables $x=e^{i \varphi}$ we arrive at

$$
\begin{equation*}
I=-\left|V_{0 n}\right|^{2} \tau \exp \left\{\frac{2\left|V_{0 n}\right|^{2} \tau}{i \zeta}\right\} \int_{0}^{2 \pi}\left(1+e^{-i \varphi}\right) \exp \left\{i \frac{2\left|V_{0 n}\right|^{2} \tau}{\zeta} \cos \varphi\right\} d \varphi \tag{3.121}
\end{equation*}
$$

which can be reduced to a standard integral representation of the Bessel functions $J_{k}(x)$ and yields

$$
\begin{equation*}
I=-2 \pi\left|V_{0 n}\right|^{2} \tau \exp \left\{\frac{2\left|V_{0 n}\right|^{2} \tau}{\zeta i}\right\}\left[J_{0}\left(\frac{2\left|V_{0 n}\right|^{2} \tau}{\zeta}\right)+i J_{1}\left(\frac{2\left|V_{0 n}\right|^{2} \tau}{\zeta}\right)\right] \tag{3.122}
\end{equation*}
$$

where the calculation has been performed assuming a positive $\tau$. The average population (3.117) now takes the form

$$
\begin{align*}
\left\langle\rho_{0}(t)\right\rangle= & \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \zeta}{2 \pi} e^{-i \zeta t} \int_{0}^{\infty} d \tau e^{i \tau \zeta} \exp \left\{-\tau \int_{0}^{\infty} g(V) 2 \pi|V|^{2}\right. \\
& \left.\exp \left\{\frac{2|V|^{2} \tau}{\zeta i}\right\}\left\{J_{0}\left(\frac{2|V|^{2} \tau}{\zeta}\right)+i J_{1}\left(\frac{2|V|^{2} \tau}{\zeta}\right)\right\} d V\right\}, \tag{3.123}
\end{align*}
$$

and implies evaluation of the integral in the exponent for a given dependence $g(V)$ of the state density on the coupling constant. If the time $t$ is short the typical frequencies $\zeta \sim 1 / t$ are big, as suggested by the uncertainty principle, and the combination $2|V|^{2} \tau / \zeta$ in the arguments under the integral in the exponent can be set to 0 , whence the integral yields the mean decay rate $W=2 \pi \int_{0}^{\infty} g(V)|V|^{2} d V$. The integration over $d \tau$ yields $1 /(\zeta+i W)$, which after the integration over $d \zeta$ results in the exponentially decaying population $\left\langle\rho_{0}(t)\right\rangle=e^{-W t}$. Now the physical meaning of the exponential and Bessel functions becomes clear: they allow "in average" for the population recurrences from the band to the state $|0\rangle$ that modify strongly the exponential decay process.

### 3.4.3 The Long-Time Asymptotic

We now consider an asymptotic expression for the population at long times. Note that the main contribution to the integrals comes from the domain where the arguments of the exponents are of the order of unity. For a long time $t$ the slow modes of frequencies $\zeta \sim 1 / t$ are important, which means that in (3.123) the main contribution to the integrals come from the domain $\tau \sim 1 / \zeta$ and $V \sim \zeta \sim 1 / t$. One can therefore replace $g(V)$ by its asymptotic form as $V \rightarrow 0$. This is consistent with the fact that at the long-time asymptotic the main part of the population is located at the band levels weakly bound to the state $|0\rangle$ which are present in abundance.

Let us consider the case $g(V) \simeq g_{\alpha} V^{-\alpha}$ for small $V$. Note that the power index $\alpha$ should be within certain limits $1 \leq \alpha<3$, since the state density $g(V)$ is supposed to satisfy two main conditions: (i) the mean transition rate, given by the integral $W=2 \pi \int g(V) V^{2} d V$ must be a finite number, which implies that the integral $\int g_{\alpha} V^{-\alpha} V^{2} d V$ converges at the lower limit $V \rightarrow 0$, and (ii) the total state density of weakly coupled states given by the integral $\int g(V) d V$ is infinite, which implies that as $V \rightarrow 0$ the integral $\int g_{\alpha} V^{-\alpha} d V$ diverges. On the contrary, the integral $\int g_{\alpha} V^{-\alpha} V d V$ of (3.109), which gives the expectation value of the number of resonant levels, may diverge or converge, thus yielding two qualitatively different knds of asymptotic behavior of the population $\left\langle\rho_{0}(t)\right\rangle$.

We also note that for some problems the density of states can indeed be given exactly by a power function of the coupling constant, as is the case for the cold Rydberg atoms with $g(V) \sim V^{-2}$ dependence discussed in Sect2.5.5. Such a distribution yields divergence of integrals $\int g(V) V^{2} d V$ and $\int g(V) V d V$ at the upper limit $V \rightarrow \infty$ and implies a cut-off procedure for strong couplings. In other problems the power dependence is valid just for the asymptotic at $V \rightarrow 0$. It turns out that convergence of the integrals is often important for physical conclusions, and therefore the replacement of the density $g(V)$ by its asymptotic dependence has to be done in a way, which excludes these artificial divergencies. One of possible ways is to take $g(V)=g_{\alpha} V^{-\alpha} e^{-(V / h)^{2}}$, where $h=\left[W / \pi g_{\alpha} \Gamma\left(\frac{3-\alpha}{2}\right)\right]^{2 /(\alpha-3)}$ is the cut-off parameter consistent with the notation $W=2 \pi \int g(V) V^{2} d V$, which can be set to 1 by a proper choice of time units, and $\Gamma(X)$ is the gamma function. This allows us to calculate the integral in the exponent of (3.123):

$$
\begin{equation*}
\int_{0}^{\infty} g(V)|V|^{2} \exp \left\{\frac{2|V|^{2} \tau}{\zeta i}\right\}\left[J_{0}\left(\frac{2|V|^{2} \tau}{\zeta}\right)+i J_{1}\left(\frac{2|V|^{2} \tau}{\zeta}\right)\right] d V \tag{3.124}
\end{equation*}
$$

analytically in terms of hypergeometric functions, and obtain the asymptotic expressions for large $\tau / \zeta \sim t^{2}$ directly from there. However, we present here an approximate albeit physically meaningful approach, which demonstrates directly the nature of the quantities that govern the long-time behavior of the system.

The replacement $\zeta \rightarrow z / t, \tau \rightarrow \theta t$ in (3.123) yields

$$
\begin{align*}
\left\langle\rho_{0}(t)\right\rangle= & \int_{-\infty+i \nu}^{\infty+i \nu} \int_{0}^{\infty} \frac{d z d \theta}{2 \pi} \exp \left\{-2 \pi t \theta \int_{0}^{\infty} g(V)|V|^{2} e^{-i|V|^{2} v}\right. \\
& {\left.\left[J_{0}\left(|V|^{2} v\right)+i J_{1}\left(|V|^{2} v\right)\right] d V-i z+i \theta z\right\} } \tag{3.125}
\end{align*}
$$

where $v=2 \theta t^{2} / z$ is a large value for long $t$. We note that the real part of the integral in the exponent is an even function of $v$, whereas the imaginary part is an odd one, which makes the entire expression real, as the
population should be. One could employ here the asymptotic form of Bessel functions $J_{n}(x) \simeq \sqrt{2 / \pi x} \cos (x-n \pi / 2-\pi / 4)$ unless small values of $V$ were important. We therefore divide the integration interval into two parts $\int_{0}^{\infty}=\int_{0}^{1 / \sqrt{v}}+\int_{1 / \sqrt{v}}^{\infty}$. The first integral corresponds to small values of $|V|^{2} v$ and we set this combination to zero in the integrand, whereas the second integral corresponds to $|V|^{2} v \gg 1$ which allows one to make use of the aforesaid asymptotic form of the Bessel functions. This results in

$$
\begin{equation*}
\exp \left(-i|V|^{2} v\right)\left[J_{0}\left(|V|^{2} v\right)+i J_{1}\left(|V|^{2} v\right)\right]=e^{i \pi / 4} \sqrt{2 / \pi|V|^{2} v} \tag{3.126}
\end{equation*}
$$

and after substitution to (3.125) yields

$$
\begin{align*}
\left\langle\rho_{0}(t)\right\rangle= & \int_{-\infty+i \nu}^{\infty+i \nu} \int_{0}^{\infty} \frac{d z d \theta}{2 \pi} \exp \left\{-t \theta \int_{0}^{1 / \sqrt{v}} 2 \pi g(V)|V|^{2} d V\right. \\
& \left.-e^{i \pi / 4} \sqrt{\theta z / \pi} \int_{1 / \sqrt{v}}^{\infty} 2 \pi g(V)|V| d V-i z+i \theta z\right\} \tag{3.127}
\end{align*}
$$

In the integrals over $d V$ one now recognizes two meaningful quantities. The first integral gives the total decay rate to the states that have coupling matrix elements $V$ less than $1 / \sqrt{v}$, that is less than or of the order of the inverse time $1 / t$. For these states $V t<1$, and hence they are not populated but work just as an empty reservoir for the exponential decay of the state $|0\rangle$. At large $v$ this integral is small, since the mean decay rate $W$ converges at the lower limit, and for $2 \pi g(V) \simeq g_{\alpha} V^{-\alpha}$ it amounts to $g_{\alpha} v^{(\alpha-3) / 2} /(3-\alpha)$. The second integral gives the expectation value of the number of resonances with the coupling $V>1 / \sqrt{v}$, calculated with the help of the two-level criterion (3.109). We have to consider two physically different cases, $\alpha<2$ when it converges as $V \rightarrow 0$, and $\alpha>2$ when it diverges. For the first case the integral approaches its asymptotic value $P_{r}=2 \pi \int_{0}^{\infty} g(V) V d V$, with the deviation $\int_{1 / \sqrt{v}}^{0} 2 \pi g(V)|V| d V \simeq g_{\alpha} v^{(\alpha-2) / 2} /(\alpha-2)$, and for the second case $\int_{1 / \sqrt{v}}^{\infty} 2 \pi g(V)|V| d V=$ const $+g_{\alpha} v^{(\alpha-2) / 2} /(\alpha-2)$, and it tends to infinity like the power of time $t^{\alpha-2}$. Later on we consider separately the important case $\alpha=2$. We take into account the specific form of the integrals and write (3.127) in the form

$$
\begin{align*}
\left\langle\rho_{0}(t)\right\rangle= & \int_{-\infty+i \nu}^{\infty+i \nu} \int_{0}^{\infty} \frac{d z d \theta}{2 \pi} \exp \left\{-i z+i \theta z-e^{i \pi / 4} \sqrt{\theta z / \pi} P_{r}\right. \\
& \left.-z^{\frac{3-\alpha}{2}} \theta^{\frac{\alpha-1}{2}} t^{\alpha-2} A-\theta^{\frac{\alpha-1}{2}} z^{\frac{3-\alpha}{2}} t^{(\alpha-2)} B\right\} \tag{3.128}
\end{align*}
$$

where $A=g_{\alpha} 2^{(\alpha-3) / 2} /(3-\alpha), B=e^{-i \pi / 4} g_{\alpha} 2^{(\alpha-2) / 2} / \sqrt{\pi}(\alpha-2)$, and $P_{r}$ either is the expectation number of resonances for $\alpha<2$ or denotes an integration constant for $\alpha>2$.

## Incomplete Decay

We now evaluate the integrals (3.128) taking into account the parity of the exact expression (3.125) with respect to $v \sim \theta / z$ by considering only the real part of the integrals. For $\alpha<2$ we expand the exponent of the integral (3.128) in powers of $t^{\alpha-2}$ and arrive at

$$
\begin{align*}
\left\langle\rho_{0}(t)\right\rangle= & \operatorname{Re}\left\{\int_{-\infty+i \nu}^{\infty+i \nu} \int_{0}^{\infty} \frac{d z d \theta}{2 \pi} e^{-i z+i \theta z-\sqrt{-i \theta z / \pi} P_{r}}\right. \\
& {\left.\left[1+z^{\frac{3-\alpha}{2}} \theta^{\frac{\alpha-1}{2}} t^{\alpha-2} C_{\alpha}+\ldots\right]\right\}, } \tag{3.129}
\end{align*}
$$

where $C_{\alpha}=g_{\alpha} 2^{(\alpha-3) / 2}\left[\sqrt{2 / \pi} e^{-i \pi / 4} /(\alpha-2)-1 /(3-\alpha)\right]$ is a numerical constant. This shows that the population $\left\langle\rho_{0}(t)\right\rangle$ of the level approaches its asymptotic value

$$
\begin{equation*}
\left\langle\rho_{0}(\infty)\right\rangle=\operatorname{Re}\left[\int_{-\infty+i \nu}^{\infty+i \nu} \int_{0}^{\infty} \frac{d z d \theta}{2 \pi} e^{-i z+i \theta z-\sqrt{-i \theta z / \pi} P_{r}}\right] \tag{3.130}
\end{equation*}
$$

as an inverse power $t^{\alpha-2}$ of time. After the replacement $\theta \rightarrow-x / i z$ we perform integration over $d z$ and arrive at

$$
\begin{equation*}
\left\langle\rho_{0}(\infty)\right\rangle=\operatorname{Re}\left[\int_{0}^{\infty} d x e^{-x-\sqrt{x / \pi} P_{r}}\right] \simeq \frac{2 \pi}{P_{r}^{2}} \tag{3.131}
\end{equation*}
$$

where the last equality is the asymptotic form for large $P_{r}$ of the exact expression

$$
\begin{equation*}
\left\langle\rho_{0}(\infty)\right\rangle=1+\left[\operatorname{Erf}\left(\frac{P_{r}}{2 \sqrt{\pi}}\right)-1\right] \frac{P_{r}}{2} \exp \left(\frac{P_{r}^{2}}{4 \pi}\right) \tag{3.132}
\end{equation*}
$$

depicted in Fig. 3.22(a). Here $\operatorname{Erf}(x)$ is the error function.
Let us consider (3.131) in the context of the discussion of the number of levels in resonance Sect.3.2.2. For a uniform spectrum and interaction the total number $\mathcal{N}$ of indeed resonant levels is not given by the expectation value $P_{r} \sim V g$ calculated with the help of the two-level resonance criterion, but amounts to the square of this number $\mathcal{N} \sim(V g)^{2}$. Therefore for large $P_{r}$ in the case of a uniform population distribution among the resonance levels each of them, including the state $|0\rangle$, has on average $\langle\rho(t)\rangle \sim 1 / \mathcal{N} \sim 1 / P_{r}^{2}$. Equation (3.131) shows that it is also the case for a random spectrum with the total number of resonant levels given by the square of the expectation number of two-level resonances $\mathcal{N}=P_{r}^{2} / 2 \pi$. For a small probability of resonance $P_{r} \ll 1$ only a small part of the systems comprising the ensemble


Fig. 3.22. (a) Dependence of the asymptotic population on the expectation number $P_{r}$ of two-level resonances. (b) Dependence of the asymptotic value of the population on the density of states for the case $2 \pi g(V)=g_{2} V^{-2}$ (solid line) and for the case $2 \pi g(V)=g_{2} V^{-2} e^{-V^{2}}$ (dotted line). (c) Dependence of the power index $\beta$ of the asymptotic decay law $\left\langle\rho_{0}(t)-\rho_{0}(\infty)\right\rangle \sim t^{-\beta}$ on the power index $\alpha$ of the asymptotic dependence $2 \pi g(V) \simeq g_{\alpha} V^{-\alpha}$ of the state density on the coupling $V$ at the limit $V \rightarrow 0$.
satisfies the two-level resonance condition $V>\Delta$, which implies that each of the systems has at most one resonance whence the population of the state $|0\rangle$ is of the order of $1 / 2$. The two-level criterion valid for this case yields $\left\langle\rho_{0}(\infty)\right\rangle=1-P_{r} / 2$ in agreement with (3.132), since $\operatorname{Erf}(0)=0$.

## Complete Decay

For $\alpha>2$ the integral $\int g(V) V d V$ diverges, which implies that an infinite number of states are in resonance with the level. In consequence the decay of the level should be complete. In order to show this we keep in the exponent of (3.96) only the leading term $\sim t^{\alpha-2}$ and arrive at
$\left\langle\rho_{0}(t)\right\rangle=\operatorname{Re}\left\{\int_{-\infty+i \nu}^{\infty+i \nu} \int_{0}^{\infty} \frac{d z d \theta}{2 \pi} e^{-i z+i \theta z-\sqrt{-i \theta z / \pi} P_{r}} \exp \left[z^{\frac{3-\alpha}{2}} \theta^{\frac{\alpha-1}{2}} t^{\alpha-2} C_{\alpha}\right]\right\}$,
where $C_{\alpha}$ denotes the same numerical constant as in (3.129). We note that the main contribution to the integral comes from the domain $z \sim 1$, $\theta \sim t^{-2(\alpha-2) /(\alpha-1)}$, and therefore two terms in the exponent containing the combination $\theta z$ can be omitted, being small at long $t$. This yields

$$
\begin{equation*}
\left\langle\rho_{0}(t)\right\rangle=\operatorname{Re}\left\{\int_{-\infty+i \nu}^{\infty+i \nu} \int_{0}^{\infty} \frac{d z d \theta}{2 \pi} e^{-i z} \exp \left[z^{\frac{3-\alpha}{2}} \theta^{\frac{\alpha-1}{2}} t^{\alpha-2} C_{\alpha}\right]\right\} \tag{3.134}
\end{equation*}
$$

which, after the replacement $\theta \rightarrow \vartheta t^{-2(\alpha-2) /(\alpha-1)}$, results in

$$
\begin{equation*}
\left\langle\rho_{0}(t)\right\rangle=t^{-2(\alpha-2) /(\alpha-1)} \operatorname{Re}\left\{\int_{-\infty+i \nu}^{\infty+i \nu} \int_{0}^{\infty} \frac{d z d \vartheta}{2 \pi} e^{-i z} \exp \left[z^{\frac{3-\alpha}{2}} \vartheta^{\frac{\alpha-1}{2}} C_{\alpha}\right]\right\} \tag{3.135}
\end{equation*}
$$

where the integral amounts to a time-independent constant. Thus the ensemble average population decays to zero as the power $2(\alpha-2) /(\alpha-1)$ of time, that is $\left\langle\rho_{0}(t)\right\rangle \sim t^{-2(\alpha-2) /(\alpha-1)}$ as shown in Fig. 3.21.

We note that in both cases $\alpha>2$ and $\alpha<2$ the asymptotic time dependence of the population becomes less and less rapid when $\alpha$ approaches the value of 2 , although the asymptotic population $\left\langle\rho_{0}(\infty)\right\rangle$ changes abruptly from a finite value to 0 . This means that the particular case $\alpha=2$ when the expectation value of resonances $\int g(V) V d V$ diverges, whereas the time dependence vanishes, requires special consideration.

## The Special Case $\alpha=2$

In the case $\alpha=2$ the situation is more difficult. Efforts are required for the determination of the long-time asymptotic $\left\langle\rho_{0}(t)-\rho_{0}(\infty)\right\rangle$ which shows how the population $\left\langle\rho_{0}(t)\right\rangle$ approaches the limiting value, although it is relatively easy to find the asymptotic value $\left\langle\rho_{0}(\infty)\right\rangle$ itself. For $2 \pi g(V)=g_{2} V^{-2}$ the integral

$$
\begin{equation*}
I=2 \pi \int_{0}^{\infty} g(V)|V|^{2} \exp \left(-i|V|^{2} v\right)\left[J_{0}\left(|V|^{2} v\right)+i J_{1}\left(|V|^{2} v\right)\right] d V \tag{3.136}
\end{equation*}
$$

entering (3.125) should not be split into two parts as we have done in (3.127), since it has very peculiar behavior. Indeed, after the substitution $2 \pi g(V)=$ $g_{2} V^{-2}$ and the replacement $V \rightarrow x \sqrt{v}(3.136)$ yields

$$
\begin{align*}
I & =g_{2} \sqrt{1 / v} \int_{0}^{\infty} \exp \left(\mp i x^{2}\right)\left[J_{0}\left(x^{2}\right) \pm i J_{1}\left(x^{2}\right)\right] d x \\
& =g_{2} \sqrt{z / \theta t^{2} \pi} e^{\mp \pi / 4}\left(-2+C_{E} \pm \frac{i}{2} \pi+5 \ln (2)\right) \tag{3.137}
\end{align*}
$$

where the upper sign corresponds to positive $z$ in the expression $v=2 \theta t^{2} / z$, and $C_{E} \simeq 0.577$ is the Euler constant. After substitution of (3.137) into (3.125) the time dependence disappears and with the allowance of the symmetry we arrive at

$$
\begin{align*}
\left\langle\rho_{0}(t)\right\rangle= & \operatorname{Re} \int_{0}^{\infty} \int_{0}^{\infty} \frac{d z d \theta}{\pi} \exp \{-i z+i \theta z- \\
& \left.g_{2} \sqrt{-i z \theta / \pi}\left(-2+C_{E}+\frac{i}{2} \pi+5 \ln (2)\right)\right\} \tag{3.138}
\end{align*}
$$

which after the replacement $\theta \rightarrow i \theta / z$ and integration of $\sin z / \pi z$ over $d z$ results in

$$
\begin{equation*}
\left\langle\rho_{0}(\infty)\right\rangle=\operatorname{Re} \int_{0}^{\infty} d \theta \exp \left\{-\theta-g_{2} \sqrt{\frac{\theta}{\pi}}\left(-2+C_{E}+\frac{i}{2} \pi+5 \ln (2)\right)\right\} \tag{3.139}
\end{equation*}
$$

This integral can be expressed in terms of the error function similar to (3.132) which now has a complex-valued argument, and in Fig. 3.22 we show the asymptotic population (3.139) as a function of the constant $g_{2}$. One sees that at $g_{2}=0$ the population remains on the level $\left\langle\rho_{0}(\infty)\right\rangle=1$, as it should be if there are no levels present in the band, whereas for larger densities it decreases as a power $\left\langle\rho_{0}(\infty)\right\rangle \sim g_{2}^{-2}$.

However, one has to take this result with caution, since it does not yield any time dependence of the population, or in other words, for such a situation the asymptotic distribution is already attained at $t=0$. This is an artefact of the model $2 \pi g(V)=g_{2} V^{-2}$, which yields the expectation number of resonances $\int g(V) V d V$ logarithmically diverging at both upper and lower limits. Thus the strong and the weak couplings $V$ are equally important, and in some way they compensate for each other in the population dynamics. Therefore, in order to determine the asymptotic time dependence of the level population $\left\langle\rho_{0}(t)-\rho_{0}(\infty)\right\rangle$ we have to consider the integral

$$
\begin{equation*}
I=\int_{0}^{\infty} e^{-V^{2}} \exp \left(-i|V|^{2} v\right)\left[J_{0}\left(|V|^{2} v\right)+i J_{1}\left(|V|^{2} v\right)\right] d V \tag{3.140}
\end{equation*}
$$

where by introducing the cut-off factor $e^{-V^{2}}$ we have taken into account that the coupling $V$ cannot have an infinite size. In this way we separate the artificial influence of strongly coupled states $(V \rightarrow \infty)$ on the long time behavior of $\left\langle\rho_{0}(t)\right\rangle$, which for the case $\alpha=2$ by accident has compensated completely the physically important contribution of the weakly coupled states.

Actually now, in the case $\alpha=2$ under consideration, even the representation (3.140), which was convenient for the splitting of the integral into two parts given in (3.127), requires rather involved analysis. The reason of this complexity becomes clear when we return to an earlier stage of the analysis and consider our initial representation (3.119), with $2 \pi g(V)=g_{2} V^{-2} e^{-V^{2}}$, which reads

$$
\begin{equation*}
\int_{0}^{\infty} \frac{g_{2}}{V^{2}} e^{-V^{2}} \int_{C} \frac{d x}{2 \pi} \exp \left\{\frac{4 i V^{2} \zeta \tau}{4 x^{2}-\zeta^{2}}\right\} d V \tag{3.141}
\end{equation*}
$$

Earlier, in order to get rid of the diverging singularity $V^{-\alpha}$ of the state density (with $\alpha=2$ in (3.141)), we have performed, although not explaining the reason, the integration by parts over $d x$ along the contour $C$ :

$$
\begin{equation*}
\int_{0}^{\infty} d V \frac{g_{2}}{V^{2}} e^{-V^{2}}\left(\left.\frac{x}{2 \pi} e^{\frac{4 i V^{2} \zeta \tau}{4 x^{2}-\zeta^{2}}}\right|_{\text {beg. } C} ^{\text {end. } C}-\int_{\text {beg. } C}^{\text {end. } C} d x \frac{x}{2 \pi} \frac{d}{d x} e^{\frac{4 i V^{2} \zeta \tau}{4 x^{2}-\zeta^{2}}}\right) \tag{3.142}
\end{equation*}
$$

As long as the integrand is an analytical function of the variable $x$ the beginning and the end of the contour $C$ correspond to the same point of the complex plane and the first term in parentheses vanishes whereas the second term becomes an integral along a closed loop. But if in the second term we try to change the order of integration and calculate the integral over $d V$ first, we obtain a non-analytic function

$$
\begin{align*}
I & =\frac{g_{2}}{\pi} \int_{\text {beg. } C}^{\text {end. } C} d x \frac{32 i \zeta \tau x^{2}}{\left(4 x^{2}-\zeta^{2}\right)^{2}} \int_{0}^{\infty} d V e^{\frac{4 i V^{2} \zeta \tau}{4 x^{2}-\zeta^{2}}-V^{2}} \\
& =\frac{g_{2}}{\sqrt{\pi}} \int_{\text {beg. } C}^{\text {end. } C} d x \frac{32 i \zeta \tau x^{2}}{\left(4 x^{2}-\zeta^{2}\right)^{3 / 2} \sqrt{4 i \zeta \tau-4 x^{2}+\zeta^{2}}}, \tag{3.143}
\end{align*}
$$

which does not return to its initial value after circumvention around the point $x=\zeta / 2$ but acquires a phase factor $e^{3 i \pi}$. This example shows the typical difficulty which arises when the ensemble average technique is employed the analytical properties of the functions before and after the average may be different, and hence the change of the order of integration might be illegal. The simplest way to avoid the problem in the case under consideration is to choose a different integration contour. The simplest suitable contour $C$ circumvents both singularities at $x= \pm \zeta / 2$ in the opposite directions andthus has a figure-eight-like shape, as shown in Fig. 3.23.


Fig. 3.23. Integration contour $C$ suitable for ensemble average. The averaged function returns to the initial value after circumventing both singularities of (3.125) in the complex plane (a) of the variable $u$. After the replacement $u \rightarrow \sqrt{y}$ this contour coincides with the Pochhammer contour (b), well-known in the theory of hypergeometric functions.

By scaling the variable $x=u \zeta / 2$ we change the size of $C$ such that it goes around the points $u= \pm 1$, and after the replacement $u=\sqrt{y}$ this contour takes the Pochhammer shape shown in Fig. 3.23. The integral itself coincides with the integral representation of the hypergeometric function ${ }_{2} F_{1}(a, b ; c ; X)$

$$
\begin{equation*}
\int_{C} \frac{\sqrt{y}(y-1)^{-3 / 2}}{\sqrt{1-y(1+4 i \tau / \zeta)^{-1}}} d y=2 \pi i{ }_{2} F_{1}\left(\frac{1}{2}, \frac{3}{2} ; 1 ; \frac{1}{1+4 i \tau / \zeta}\right) \tag{3.144}
\end{equation*}
$$

and can also be given in different forms either as a hypergeometric function of a different argument or as an elliptic function $E(X)$, which reads

$$
\begin{align*}
I & =\frac{-2 g_{2} \sqrt{\pi} \tau}{\sqrt{1+4 i \tau / \zeta}}{ }_{2} F_{1}\left(\frac{1}{2}, \frac{3}{2} ; 1 ; \frac{1}{1+4 i \tau / \zeta}\right) \\
& =i g_{2} \sqrt{\pi} \zeta \sqrt{\frac{i \tau}{\zeta}}{ }_{2} F_{1}\left(-\frac{1}{2}, \frac{1}{2} ; 1 ;-\frac{i \zeta}{4 \tau}\right) \\
& =\frac{2 i g_{2}}{\sqrt{\pi}} \zeta \sqrt{\frac{i \tau}{\zeta}} E\left(-\frac{i \zeta}{4 \tau}\right) \tag{3.145}
\end{align*}
$$

Note that when discussing (3.124) earlier we mentioned that for $\alpha \neq 0$ the integral $I$ can also be given in terms of hypergeometric functions, and indeed, by considering (3.143) for this case one obtains a result similar to (3.145), which however cannot be reduced to the elliptic function.

Substitution of (3.145) to (3.117) yields

$$
\begin{equation*}
\left\langle\rho_{1}(t)\right\rangle=\int_{-\infty+i \nu}^{\infty+i \nu} \int_{0}^{\infty} \frac{d \tau d \zeta}{2 \pi} \exp \left\{-i \zeta t+i \tau \zeta-i \zeta \frac{g_{2}}{\sqrt{\pi}} \sqrt{\frac{4 \tau}{-i \zeta}} E\left(-\frac{i \zeta}{4 \tau}\right)\right\} \tag{3.146}
\end{equation*}
$$

and after the replacement $\zeta \rightarrow-2 i z \sqrt{\theta} ; \tau \rightarrow-\sqrt{\theta} / 2 z$ results in

$$
\begin{equation*}
\left\langle\rho_{1}(t)\right\rangle=\int_{-i \infty-\nu}^{i \infty-\nu} \int_{0}^{\infty} \frac{d \theta d z}{2 \pi z} \exp \left\{-2 z \sqrt{\theta} t-\theta-\frac{2 g_{2}}{\sqrt{\pi}} \sqrt{\theta} E\left(z^{2}\right)\right\} \tag{3.147}
\end{equation*}
$$

One can check that such a replacement does not affect the convergence of the integrals.

In Fig. 3.24 we show the absolute value of the integrand (3.147), and the phase of the hypergeometric function ${ }_{2} F_{1}\left(-1 / 2,1 / 2 ; 1 ; z^{2}\right)=2 E\left(z^{2}\right) / \pi$ in the complex plane of the variable $z$. One clearly sees that the initial contour $C$ can be moved toward large $z$ and can be transformed to a contour $C^{\prime}$ consisting of two parts. The first part is a loop around the pole at $z=0$. This yields a time-independent contribution corresponding to the asymptotic population

$$
\begin{equation*}
\left\langle\rho_{0}(\infty)\right\rangle=\int_{0}^{\infty} d \theta \exp \left\{-\theta-\frac{g_{2}}{\sqrt{\pi}} \sqrt{\theta}\right\} \tag{3.148}
\end{equation*}
$$

which can be given in terms of the error function by analogy to (3.132), and which is similar in all the main features to the population shown in Fig. 3.22. The second part of the contour $C^{\prime}$ goes around the branching point of the hypergeometric function at $z=1$ and allows for the deviation $\left\langle\rho_{0}(t)-\rho_{0}(\infty)\right\rangle$ decreasing in time. Since the hypergeometric function at different sides of the real axis takes complex conjugated values, this contour integral can be transformed to a definite integral along the real axis from $z=1$ to $z=$ $\infty$ provided we keep only the imaginary part of the integrand. For long $t$ the integral can be evaluated approximately. Indeed, since the limit $t \rightarrow \infty$ corresponds to $\sqrt{\theta} \rightarrow 0$ the term $\theta$ in the exponent of the integral


Fig. 3.24. The structure of the integrand (3.147) in the complex plane $z$. (a) The absolute value of the integrand shows that the initial integration contour $C=$ $(-i \infty-\nu, i \infty-\nu)$ can be transformed to the contour $C^{\prime}$, which consists of two parts and comprises a loop around the pole at $z=0$ that gives the asymptotic value of $\left\langle\rho_{0}(t)\right\rangle$ and a contour that goes around the branching point $z=1$ shown in (b) which yields the deviation $\left\langle\rho_{0}(t)-\rho_{0}(\infty)\right\rangle$ vanishing at long $t$. (b) The phase of the hypergeometric function of (3.147) is discontinuous on the real axis, starting from the branching point at $z=1$.

$$
\begin{equation*}
\left\langle\rho_{0}(t)-\rho_{0}(\infty)\right\rangle=-2 \operatorname{Im} \int_{1}^{\infty} \int_{0}^{\infty} \frac{d \theta d z}{2 \pi z} e^{-2 z \sqrt{\theta} t-\theta-\frac{2 g_{2}}{\sqrt{\pi}} \sqrt{\theta} E\left(z^{2}\right)} \tag{3.149}
\end{equation*}
$$

can be omitted as small, and the integration over $d \theta$ yields

$$
\begin{equation*}
\left\langle\rho_{0}(t)-\rho_{0}(\infty)\right\rangle \simeq-2 \operatorname{Im} \int_{1}^{\infty} \frac{d z}{\pi z\left(2 z t+\frac{2 g_{2}}{\sqrt{\pi}} E\left(z^{2}\right)\right)^{2}} \tag{3.150}
\end{equation*}
$$

The integrand tends to a real number $1 / 4 \pi z^{3} t^{2}$ for long $t$, and the first non-vanishing correction $g_{2} \operatorname{Im}\left[E\left(z^{2}\right)\right] / 2 \pi \sqrt{\pi} z^{4} t^{3}$ can be determined by the Taylor expansion over $g_{2}$ and, after integration, yields

$$
\begin{equation*}
\left\langle\rho_{0}(t)-\rho_{0}(\infty)\right\rangle \simeq \frac{g_{2}}{\pi \sqrt{\pi} t^{3}} \int_{1}^{\infty} \operatorname{Im}\left[E\left(z^{2}\right)\right] \frac{d z}{z^{4}}=\frac{2 g_{2}}{9 \pi \sqrt{\pi} t^{3}} \tag{3.151}
\end{equation*}
$$

One sees, in spite of the fact that for the case $\alpha=2$ the expectation number of two-level resonances $P_{r}$ is infinite, the decay of the level $|0\rangle$ is not complete, although the population of this level approaches its finite asymptotic value $2 \pi / g_{2}^{2}$ as $t^{-3}$, which is much faster than for any other $\alpha$, as shown in Fig. 3.22. It worth mentioning that here the coefficient $g_{2}$ plays the role of $P_{r}$ for other $\alpha$.

## 4 Two-Band System

Thus far we have been considering the level-band system, for which an explicit analytic expression (3.26) for the Fourier transform of the probability amplitude exists. However, as we have seen, even for such a system, the complexity of the problem does not allow one to directly perform the inverse Fourier transformation in the generic case, and therefore in order to get physically consistent analytic results, one needs to apply a statistical description based on the idea of ensemble averages. The situation becomes more difficult still when we consider the more complex quantum system of two coupled bands. Not only the inverse Fourier transformation, but even the derivation of an explicit, exact expression for the Fourier transforms like (3.26) turns out to be problematic. Therefore, employment of the ensemble average approach becomes a vital part of the analytic approach to this problem.

We start with the consideration of the all-order perturbation series and their graphical representation in the spirit of Feynman diagrams. However, being topologically equivalent to the standard diagrams of many-body theory where the unperturbed motion of the particle is depicted as straight lines and the perturbations as vortices with the incoming lines of the perturbation, our diagrams more closely resemble Grotrian ones, with interaction shown by lines of transitions among the points denoting the unperturbed energy levels. We believe that such a representation, performed in the spirit suggested by P.W. Anderson in his seminal paper on dynamical localization, is better suited to the case of point (0-dimensional) systems like atoms and molecules, where the free evolution does not imply propagation along a straight trajectory.

We discuss the renormalization associated with the random perturbation, and consider several particular cases of non-degenerate bands, two degenerate levels, and a degenerate level coupled to a non-degenerate band. We consider the dynamics of the total populations of the bands and the degenerate levels as well as the population distribution over the band, and demonstrate the interplay between the exponential and the coherent damping that yields oscillations in the time dependencies of the total populations and oscillations in the distribution of the population over the bands.

We dwell on the role of the correlations of the perturbation that manifest themselves in the fact that the rank of the perturbation matrix becomes much smaller relative to its order, and elucidate the role of correlations between the
matrix elements of the perturbation and the energy position of the levels to which they belong. We show how these correlations affect the time evolution of the total populations and the distribution of the populations over the bands. We conclude by considering the situation where all of the matrix elements of the perturbation have Gaussian statistics but with an individual width for each transition.

### 4.1 General Consideration

Under certain conditions, the calculation of the population dynamics in different complex quantum systems can be carried out in the same manner. We first express the Fourier transforms of the probability amplitudes to be at different quantum states in terms of the all-order perturbation series over the coupling. We multiply these series by the similar series for the complex conjugate amplitudes in order to obtain the probabilities. We then perform the ensemble average over the sizes of the couplings and thereby eliminate the majority of terms in the perturbation series. This often allows one to perform an analytic summation of the series for probabilities. We note that after the ensemble average, the particular values of the interaction matrix elements are no longer important for certain regimes, whereas the topological structure of the perturbation series, that is the number of possible sequences of adjunct transitions induced by the perturbation, starts to play a leading role.

### 4.1.1 Series and Diagrams for the Level-Band Problem

The casting of the probability amplitudes in the all-order perturbation series is the first step of the analysis, and in order to gain a deeper insight into the approach we start with the example of the level-band system, which has already been considered in the previous chapter, but this time with the help of a different technique.

## Series and Diagrams for the Probability Amplitudes

One can cast the exact expression (3.26) for the Fourier transform of the probability amplitude tof remaining in the initial state $|0\rangle$ in a power series over the coupling

$$
\begin{align*}
\psi_{0}(\varepsilon)= & \frac{i}{\varepsilon-\sum_{n} \frac{V_{0}^{n} V_{n}^{0}}{\varepsilon-\Delta_{n}}}=i\left[\frac{1}{\varepsilon}+\frac{1}{\varepsilon}\left(\sum_{n} \frac{V_{0}^{n} V_{n}^{0}}{\varepsilon-\Delta_{n}}\right) \frac{1}{\varepsilon}\right. \\
& +\frac{1}{\varepsilon}\left(\sum_{n} \frac{V_{0}^{n} V_{n}^{0}}{\varepsilon-\Delta_{n}}\right) \frac{1}{\varepsilon}\left(\sum_{n} \frac{V_{0}^{n} V_{n}^{0}}{\varepsilon-\Delta_{n}}\right) \frac{1}{\varepsilon}+\ldots \\
& \left.+\frac{1}{\varepsilon^{k+1}}\left(\sum_{n} \frac{V_{0}^{n} V_{n}^{0}}{\varepsilon-\Delta_{n}}\right)^{k}+\ldots . .\right]  \tag{4.1}\\
= & i\left[\frac{1}{\varepsilon}+\sum_{n=1}^{N} \frac{1}{\varepsilon} \frac{V_{0}^{n} V_{n}^{0}}{\left(\varepsilon-\Delta_{n}\right)} \frac{1}{\varepsilon}+\frac{1}{\varepsilon} \sum_{n, m} \frac{V_{0}^{n} V_{n}^{0}}{\left(\varepsilon-\Delta_{n}\right)} \frac{1}{\varepsilon} \frac{V_{0}^{m} V_{m}^{0}}{\left(\varepsilon-\Delta_{m}\right)} \frac{1}{\varepsilon}\right. \\
& \left.+\frac{1}{\varepsilon} \sum_{n, m, k} \frac{V_{0}^{n} V_{n}^{0}}{\left(\varepsilon-\Delta_{n}\right)} \frac{1}{\varepsilon} \frac{V_{0}^{m} V_{m}^{0}}{\left(\varepsilon-\Delta_{m}\right)} \frac{1}{\varepsilon} \frac{V_{0}^{k} V_{k}^{0}}{\left(\varepsilon-\Delta_{k}\right)} \frac{1}{\varepsilon}+\ldots\right] . \tag{4.2}
\end{align*}
$$

One sees that the generic term of the perturbation series has the structure

$$
\begin{align*}
& \frac{1}{\varepsilon} V_{0}^{n} \frac{1}{\left(\varepsilon-\Delta_{n}\right)} V_{n}^{0} \frac{1}{\varepsilon} V_{0}^{m} \frac{1}{\left(\varepsilon-\Delta_{m}\right)} V_{m}^{0} \frac{1}{\varepsilon} V_{0}^{k} \frac{1}{\left(\varepsilon-\Delta_{k}\right)} V_{k}^{0} \frac{1}{\varepsilon} \times \ldots \\
& \times \frac{1}{\varepsilon} V_{0}^{l} \frac{1}{\left(\varepsilon-\Delta_{l}\right)} V_{l}^{0} \frac{1}{\varepsilon} V_{0}^{s} \frac{1}{\left(\varepsilon-\Delta_{k}\right)} V_{s}^{0} \frac{1}{\varepsilon} \tag{4.3}
\end{align*}
$$

which can be interpreted in physical terms as a sequence of transitions induced by the interaction $V_{0}^{n}$, and therefore the entire series represents the sum over all transition trajectories. In Fig. 4.1 we show such a trajectory for a 6 -th order term

$$
\begin{equation*}
\frac{1}{\varepsilon} V_{0}^{n} \frac{1}{\left(\varepsilon-\Delta_{n}\right)} V_{n}^{0} \frac{1}{\varepsilon} V_{0}^{m} \frac{1}{\left(\varepsilon-\Delta_{m}\right)} V_{m}^{0} \frac{1}{\varepsilon} V_{0}^{k} \frac{1}{\left(\varepsilon-\Delta_{k}\right)} V_{k}^{0} \frac{1}{\varepsilon} \tag{4.4}
\end{equation*}
$$

as an example. The system initially in the state $|0\rangle$ makes a transition to the


Fig. 4.1.
state $|n\rangle$ with the transition amplitude given by the matrix element $V_{0}^{n}$ and returns back to the state $|0\rangle$ with the amplitude $V_{n}^{0}$ in order to perform two more "round trip" transitions, to the state $|m\rangle$ and to the state $|k\rangle$ consecutively. The visiting of each level $|l\rangle$ is associated with a fraction $1 /\left(\varepsilon-\Delta_{l}\right)$ that corresponds to a pole in the complex plane of $\varepsilon$ at the level energy position $\Delta_{l}$ and yields an oscillating exponent $e^{i \Delta_{l} t}$ after the inverse Fourier transformation, whereas the energy position $\Delta_{0}$ of the level $|0\rangle$ is taken to be zero and gives the fraction $1 / \varepsilon$ corresponding to time integration.


Fig. 4.2.

Hereafter, instead of writing down the terms of the perturbation series, which becomes more and more cumbersome as the order of the perturbation increases, we will draw the corresponding diagrams, where each transition between two states $|e\rangle$ and $|s\rangle$ with the amplitude $V_{e}^{s}$ is shown by a wave line and points on the ends of these lines denote the fractions $1 /\left(\varepsilon-\Delta_{l}\right)$. In this notation, the amplitude (4.2) is given by a sum over all possible "bush-like" trajectories of any possible length that start and end at the point $|0\rangle$, as shown in Fig. 4.2.

By $X_{0}(\varepsilon)$ we denote here the sum of all diagrams that start and end with state $|0\rangle$, which coincides with the probability amplitude $\psi_{0}(\varepsilon)=X_{0}(\varepsilon)$ for the particular level-band problem under the consideration.

## Repetition of the Transitions

The representation of the probability amplitudes in the form of a diagramatic series not only provides a convenient way to handle cumbersome expressions, but occasionally also suggests a hint for practical calculations. For example, the role of repeating transitions between a pair of states can be readily understood with the help of diagrams. This relies on the comprehension that each term of the perturbation series corresponds to a sequence of transitions, and each transition requires a time $\tau$ which is of the order either of the inverse transition probability $\tau \sim 1 / W$ or of the inverse typical coupling $\tau \sim 1 / V$, depending on which is longer. Therefore at a given moment $t$ the total number of transitions that have occurred in the system is limited, and the typical length of the diagram can be estimated as $N(t) \sim t / \tau$.

If the number $N(t)$ is much less than the total number $\mathcal{N} \simeq W / \delta \sim$ $V^{2} / \delta^{2} \sim P_{r}^{2}$ of the levels involved in the process discussed earlier in Sect.3.2.2
and Sect.3.4.3 (3.132), the probability to encounter a trajectory which arrives twice or more times at the same state $|n\rangle$ of the band is small. Although these trajectories enter into the sum over all diagrams, their relative contribution is negligible for short time-scales. Therefore the ensemble average procedure yields only the mean squared matrix elements $\left\langle V_{0}^{n} V_{n}^{0}\right\rangle$, corresponding to the "forward-and-back" transition from the level $|0\rangle$ to a state $|n\rangle$, where the contribution of the higher-order averages remains unimportant. This is in complete agreement with the result of (3.123) and the discussion thereafter, showing that for the short-time regime, while the recurrences or revivals have not yet occurred, only the mean square coupling $\int g(V) V^{2} d V$ is of importance.

This reasoning brings us to the conclusion that by taking into account repeating transitions to the same states, we allow for the influence of recurrences and revivals on the population dynamics of a multilevel system, that definitely play an important role over a long period of time but do not contribute to the short-time dynamics. As we have already seen in Sect.3.3.2 the revival regime arises at times $t \sim 1 / \delta \sim g$ of the order of the inverse spacing $1 / \delta$ between the neighboring energy levels, that is the mean state density $g$. This means that one can ignore the returns for $t \ll g$ which also corresponds to the regime where one can make use of the continuous-band model.

## Diagrams for Populations

As we have already mentioned when considering the level-band system with a complicated and irregular structure of the detunings $\Delta_{n}$ and of the transition matrix elements $V_{0}^{n}$, the value of the probability amplitude $\psi_{n}(t)$ of a state $n$ may and usually does vanish after average over the ensemble. Therefore, dealing with complex systems we have to consider mainly the probabilities (3.111), that are "good" or "self-averaging" quantities, and draw the corre-


Fig. 4.3.
sponding terms of the perturbation series for the populations. We do this by superimposing the diagrams corresponding to the Fourier transform $\psi_{0}(\varepsilon)$
of the amplitude $\psi_{0}(t)$ with the diagrams for the Fourier transforms $\psi_{0}(\xi)$ of the complex conjugate amplitude $\psi_{0}^{*}(t)$ as shown in Fig. 4.3. In order to distinguish between these two types of diagrams we use undulating lines for the trajectories of $\psi_{0}(\varepsilon)$ and dashed lines for those of $\psi_{0}(\xi)$. It is clear that in a typical term of the perturbation series for populations, the undulating and the dashed trajectories have different lengths. We also note that the dashed line connecting two points, say $p$ and $q$ corresponds to the factor $\left(V_{q}^{p}\right)^{*}=V_{p}^{q}$ conjugated with the same coupling factor on the wavy line.

The technique under discussion plays no important role in the consideration of the level-band systems; it only brings to light the fact that each term of the perturbation series represents a certain trajectory of sequential jumps, and that the all-order perturbation series itself is just a sum of the amplitudes of all possible trajectories. This property is a particular case of a more general rule stating that in quantum mechanics we have to take into account the interference of all possible channels of the process. It is also a complete analog of the Feynman path integrals, where the integration is carried out in the discrete space of the energy eigenstates of the unperturbed Hamiltonian.

We note, however, that for more complex systems this approach seems to be more convenient as compared with the other methods and to yield the final results in the most straightforward and physically clear way. Indeed for the level-band system, the topologic structure of the interaction is such that a transition $0 \rightarrow n$ must be followed immediately by the back transition $n \rightarrow 0$, and therefore each matrix element $V_{0}^{n}$ entering the perturbation series together with its complex conjugate $V_{n}^{0}$ form an irreducible block $V_{0}^{n}\left(\varepsilon-\Delta_{n}\right)^{-1} V_{n}^{0}$ in the geometric series (4.2). This is no longer the case for more complex two-band systems, where the transition from a level of the first band to an energy eigenstate of the second band does not imply an immediate back transition. On the contrary, for this problem the generic trajectory travels over different states with little probability of a self-intersection occurring in the regime $t \ll g$.

### 4.1.2 Series and Diagrams for the Two-Band Problem

Let us now consider the two-band problem. We designate levels of the first band by the index $n$ and levels of the second band by index $m$, both indeces running either from $-\infty$ to $\infty$ for infinite number of states or from $-N / 2$ to $N / 2$ and from $-M / 2$ to $M / 2$ for a finite numbers of states $N+1$ and $M+1$ of the first and the second bands respectively. We assume that at $t=0$ only the state $n=0$ is populated, that is $\psi_{n}(t=0)=\delta_{n}^{0}, \psi_{m}(t=0)=0$, and by analogy to (3.10) write down the Fourier transformed Schrödinger equation

$$
\begin{align*}
& \varepsilon \psi_{n}(\varepsilon)=\Delta_{n} \psi_{n}(\varepsilon)+V_{n}^{m} \psi_{m}(\varepsilon)+i \delta_{n}^{0} \\
& \varepsilon \psi_{m}(\varepsilon)=\Delta_{m} \psi_{m}(\varepsilon)+V_{m}^{n} \psi_{n}(\varepsilon) \tag{4.5}
\end{align*}
$$

where the summation over repeating upper and lower indices is implicit, that is $V_{n}^{m} \psi_{m}=\sum_{m} V_{n}^{m} \psi_{m}$. Here, as earlier, $\Delta_{n}$ and $\Delta_{m}$ denote detunings from
the resonance of the levels of the first and second bands, respectively, and $\Delta_{n=0}$ is set to zero.

## Matrix Form

Two rectangular $(N+1) \times(M+1)$ matrices $\hat{V}$ and $\hat{V}^{+}$of an interaction with the complex conjugate elements $V_{n}^{m}=\left(\left(V^{+}\right)_{m}^{n}\right)^{*}$ together with the unperturbed Hamiltonians of the first $\hat{\Delta}_{1}$ and second $\hat{\Delta}_{2}$ bands comprise the Hamiltonian

$$
\widehat{H}=\left(\begin{array}{cc}
\hat{\Delta}_{1} & \hat{V}  \tag{4.6}\\
\hat{V}^{+} & \hat{\Delta}_{2}
\end{array}\right)
$$

of the two-band system, which can be separated into two parts: (i) the unperturbed Hamiltonian

$$
\widehat{H}_{0}=\left(\begin{array}{cc}
\hat{\Delta}_{1} & 0  \tag{4.7}\\
0 & \hat{\Delta}_{2}
\end{array}\right)
$$

with $\left(\hat{\Delta}_{1}\right)_{n}^{n^{\prime}}=\delta_{n}^{n^{\prime}} \Delta_{n}$ and $\left(\hat{\Delta}_{2}\right)_{m}^{m^{\prime}}=\delta_{m}^{m^{\prime}} \Delta_{m}$ in the unperturbed energy eigenstate representation, and (ii) the interaction

$$
\widehat{V}_{i n t}=\left(\begin{array}{cc}
0 & \hat{V}_{12}  \tag{4.8}\\
\hat{V}_{21}^{+} & 0
\end{array}\right)
$$

where $\widehat{H}_{0}$ and $\widehat{V}_{i n t}$ are both square Hermitian matrices.

## Series

In the matrix notation of (4.7)-(4.8) the Schrödinger equation (4.5) reads

$$
\begin{equation*}
\left[\varepsilon-\widehat{H}_{0}-\widehat{V}_{i n t}\right] \bar{\psi}(\varepsilon)=i \overline{\delta_{n}^{0}} \tag{4.9}
\end{equation*}
$$

where $\bar{\psi}$ denotes the vector $\left\{\psi_{-N / 2}, \ldots, \psi_{N / 2}, \psi_{-M / 2}, \ldots, \psi_{M / 2}\right\}$, and $\overline{\delta_{n}^{0}}=$ $\left\{0, \ldots, \delta_{n}^{0}, \ldots, 0, \ldots, 0\right\}$ corresponds to the initial condition. The solution of (4.9)

$$
\begin{equation*}
\bar{\psi}(\varepsilon)=\frac{i}{\varepsilon-\widehat{H}_{0}-\widehat{V}_{i n t}} \overline{\delta_{n}^{0}} \tag{4.10}
\end{equation*}
$$

given in terms of the resolvent $1 /\left(\varepsilon-\widehat{H}_{0}-\widehat{V}_{\text {int }}\right)$ means

$$
\begin{align*}
\psi_{n}(\varepsilon) & =\langle n| \frac{i}{\varepsilon-\widehat{H}_{0}-\widehat{V}_{i n t}}|n=0\rangle  \tag{4.11}\\
\psi_{m}(\varepsilon) & =\langle m| \frac{i}{\varepsilon-\widehat{H}_{0}-\widehat{V}_{i n t}}|n=0\rangle \tag{4.12}
\end{align*}
$$

Now we make use of the Taylor expansion of the resolvent operator

$$
\begin{align*}
\frac{1}{\varepsilon-\widehat{H}_{0}-\widehat{V}_{i n t}}= & \frac{1}{\varepsilon-\widehat{H}_{0}}+\frac{1}{\varepsilon-\widehat{H}_{0}} \widehat{V}_{i n t} \frac{1}{\varepsilon-\widehat{H}_{0}} \\
& +\frac{1}{\varepsilon-\widehat{H}_{0}} \widehat{V}_{i n t} \frac{1}{\varepsilon-\widehat{H}_{0}} \widehat{V}_{i n t} \frac{1}{\varepsilon-\widehat{H}_{0}}+\ldots \tag{4.13}
\end{align*}
$$

and take into account the block structure (4.6) of the Hamiltonian. The even terms of (4.13) then contribute to (4.11) and yield the probability amplitudes to find the system at the levels of the first band, whereas the odd terms correspond to (4.12), that is to the second band. This yields

$$
\begin{align*}
\psi_{n}(\varepsilon)= & \frac{i \delta_{n}^{0}}{\varepsilon-\Delta_{0}}+i \frac{1}{\varepsilon-\Delta_{0}} V_{0}^{m} \frac{1}{\varepsilon-\Delta_{m}} V_{m}^{n} \frac{1}{\varepsilon-\Delta_{n}}+\ldots \\
\psi_{m}(\varepsilon)= & i \frac{1}{\varepsilon-\Delta_{0}} V_{0}^{m} \frac{1}{\varepsilon-\Delta_{m}} \\
& +i \frac{1}{\varepsilon-\Delta_{0}} V_{0}^{m^{\prime}} \frac{1}{\varepsilon-\Delta_{m^{\prime}}} V_{m^{\prime}}^{n} \frac{1}{\varepsilon-\Delta_{n}} V_{n}^{m} \frac{1}{\varepsilon-\Delta_{m}}+\ldots, \tag{4.14}
\end{align*}
$$

where the summation over repeating indices is implicit.

## Diagrams

Comparison of (4.14) and (4.4) suggests that the amplitudes of the states $\psi_{n}(\varepsilon)$ and $\psi_{m}(\varepsilon)$ are sums of the contributions of all possible trajectories that start at the state $|0\rangle$, that is $|n=0\rangle$ of the first band and end at the states $|n\rangle$ or $|m\rangle$ respectively. In Fig. 4.4 we depict the diagrams representing the 6 -th order perturbation term for $\psi_{n}(\varepsilon)$ and the 5 -th order perturbation term for $\psi_{m}(\varepsilon)$ with the help of the same notation that has been employed for the level-band problem.


Fig. 4.4.

In contrast to the diagrams of Figs. 4.1 and 4.2 the trajectories in the two-band system do not have to return to the same level because of the interaction topology and indeed on the contrary, they travel over all possible
intermediate states. Therefore generic terms of the series (4.14) contain the majority of factors $V_{n}^{m}$ in the first power. This is also true for terms of the series for populations such as those depicted in Fig. 4.5. This has a dramatic consequence for the ensemble averaged quantities, since after the average over an ensemble of matrices $\widehat{V}_{\text {int }}$ with a symmetric, say Gaussian, distribution of the matrix elements $V_{n}^{m}$ the majority of terms in the perturbation series vanish. Indeed, if a term of the perturbation series contains even a single


Fig. 4.5.
matrix element $V_{n}^{m}$ without its complex conjugate counterpart $\left(V_{n}^{m}\right)^{*}$, it is an odd function which yields zero after the ensemble average with an even distribution. It is worth mentioning that if the mean matrix element $\left\langle V_{n}^{m}\right\rangle$ differs from zero, then the mean interaction $\left\langle\widehat{V}_{\text {int }}\right\rangle$ can be attributed to the unperturbed Hamiltonian $\widehat{H}_{0} \rightarrow \widehat{H}_{0}+\left\langle\widehat{V}_{i n t}\right\rangle$, such that the ensemble average of the rest of the perturbation $\widehat{V}_{\text {int }}-\left\langle\widehat{V}_{i n t}\right\rangle$ would yield zero for the mean values of the matrix elements.

## The Diagram Topology

We must now consider the structure of the non-vanishing terms. Generally speaking their structure is complex, once the long-time asymptotic behavior is concerned, as detailed later. Here, however, we concentrate on the regime $t \ll$ $g$ where the structure may be readily determined. Intuitively we understand that during a finite time interval only a finite number of transitions can occur in the two-band system, as was the case in the level-band system previously discussed. If the interval is short enough, the number of transitions is much less than the total number of levels involved in the population dynamics, and hence the diagram has enough space for such a location where the number of self-intersections is not larger than that needed to sustain the ensemble average. An example of a diagram of 22 -nd order for the amplitude is shown in Fig. 4.6(a). In such a diagram each transition $n \rightarrow m$ occurs together with its counterpart - the transition $m \rightarrow n$ in the opposite direction, which yields a factor $\left\langle V_{m}^{n} V_{n}^{m}\right\rangle$ after the ensemble average. There is no need to take


Fig. 4.6.
into account repetition of the transitions, since their contribution is small as long as the total number of transitions remains small compared to the total number $\mathcal{N}$ of the levels involved in the process. We illustrate this in Fig. 4.6(b) and (c) where the total number of different diagrams of type (c) is $\mathcal{N}$ times smaller than the number of the diagrams of type (b), whereas a contribution to the perturbation series of each particular 8-th order diagram of these types is of the same order of magnitude.

We now depict the diagram for a generic, non-vanishing term of the perturbation series for the amplitude without paying attention to the particular indices of the visited states. It has the tree-like structure shown in Fig. 4.7. Let us convince ourselves by directly performing the ensemble averaging that the terms of the perturbation series indeed have a tree-like structure. We first consider the matrix elements that are the closest neighbors of their conjugates, that is the parts of the perturbation series terms where the sequence of factors reads

$$
\begin{equation*}
\sum_{m} \ldots \times V_{k}^{n} \frac{1}{\left(\varepsilon-\Delta_{n}\right)} V_{n}^{m} \frac{1}{\left(\varepsilon-\Delta_{m}\right)} V_{m}^{n} \frac{1}{\left(\varepsilon-\Delta_{n}\right)} V_{n}^{l} \times \ldots \tag{4.15}
\end{equation*}
$$

The ensemble average of this part yields

$$
\begin{align*}
& \ldots \times \frac{1}{\left(\varepsilon-\Delta_{n}\right)} \sum_{m} \frac{\left\langle V_{n}^{m} V_{m}^{n}\right\rangle}{\left(\varepsilon-\Delta_{m}\right)} \frac{1}{\left(\varepsilon-\Delta_{n}\right)} \times \ldots \\
= & \ldots \times \frac{1}{\left(\varepsilon-\Delta_{n}\right)} \sum_{m} \frac{\left\langle V^{2}\right\rangle}{\left(\varepsilon-\Delta_{m}\right)} \frac{1}{\left(\varepsilon-\Delta_{n}\right)} \times \ldots, \tag{4.16}
\end{align*}
$$

where by $\left\langle V^{2}\right\rangle=\left\langle V_{n}^{m} V_{m}^{n}\right\rangle=\int\left|V_{n}^{m}\right|^{2} g\left(\left|V_{n}^{m}\right|\right) d V_{n}^{m}$ we denote the mean square coupling corresponding to the distribution function $g\left(\left|V_{n}^{m}\right|\right)$. One can


Fig. 4.7.
see that the ensemble average wipes off the matrix structure of the perturbation such that the remaining factor $\sum_{m}\left\langle V^{2}\right\rangle /\left(\varepsilon-\Delta_{m}\right)$ can be attributed to its "root" state $n$, and therefore the next two matrix elements $V_{k}^{n}$ and $V_{n}^{l}$ separated by this factor can now be considered as neighboring. We repeat this procedure again and again until all of the transition matrix elements become coupled, as shown in Fig. 4.7.

### 4.1.3 The Renormalization

The tree-like structure of the diagrams after the ensemble average in the short-time regime allows one to perform an analytical summation of the perturbation series and to reduce the problem to the solution of a set of algebraic equations for the ensemble averaged sums and average characteristics of the system, that do not contain parameters of particular systems comprising the ensemble. The key idea of the approach employs the fact that each "branch" of a "tree" has exactly the same topology as the entire "tree".

## Equations for the Averaged Amplitudes

Indeed, let us denote by $X_{m}(\varepsilon)$ the sum of all possible tree-like diagrams that start at a level $m$ and return to the level $m$ of the second band, and let us concentrate on the structure of diagrams that start at a specific level $n$ of the first band. The perturbation series for the amplitudes with the allowance of the fact that a branch also has this tree-like topology and therefore can be denoted by $X_{m}(\varepsilon)$, takes the form shown in Fig. 4.8 where we denote $\left\langle V^{2}\right\rangle$ by $V^{2}$. The first term of the series contains no transitions and yields just the


Fig. 4.8.
factor $1 /\left(\varepsilon-\Delta_{n}\right)$ corresponding to the level $n$ of the first band. The first transition starts at the level $n$, goes to a level $m$ of the second band and returns back. It yields the factor $\left\langle V^{2}\right\rangle /\left(\varepsilon-\Delta_{m}\right)\left(\varepsilon-\Delta_{n}\right)^{2}$, which can be combined with all other non-vanishing terms of the perturbation series that pass through the state $m$ and return only once to the state $n$ thus yielding the combination

$$
\begin{equation*}
\frac{\left\langle V_{n}^{m} X_{m}(\varepsilon) V_{m}^{n}\right\rangle}{\left(\varepsilon-\Delta_{n}\right)^{2}}=\frac{\left\langle V^{2}\right\rangle X_{m}(\varepsilon)}{\left(\varepsilon-\Delta_{n}\right)^{2}} \tag{4.17}
\end{equation*}
$$

Summation over $m$ yields the second term of the series shown in Fig. 4.8. The next term combines all the diagrams that return to the state $n$ twice, and thus the contribution reads

$$
\begin{align*}
& \frac{1}{\varepsilon-\Delta_{n}}\left\langle\sum_{m} V_{n}^{m} X_{m}(\varepsilon) V_{m}^{n}\right\rangle \frac{1}{\varepsilon-\Delta_{n}}\left\langle\sum_{m^{\prime}} V_{n}^{m^{\prime}} X_{m^{\prime}}(\varepsilon) V_{m^{\prime}}^{n}\right\rangle \frac{1}{\varepsilon-\Delta_{n}} \\
= & \frac{1}{\varepsilon-\Delta_{n}}\left\langle V^{2}\right\rangle \sum_{m} X_{m}(\varepsilon) \frac{1}{\varepsilon-\Delta_{n}}\left\langle V^{2}\right\rangle \sum_{m^{\prime}} X_{m^{\prime}}(\varepsilon) \frac{1}{\varepsilon-\Delta_{n}}, \tag{4.18}
\end{align*}
$$

whereas the generic term of the series corresponding to $k$ returns has the form

$$
\begin{equation*}
\frac{1}{\left(\varepsilon-\Delta_{n}\right)^{k+1}}\left\langle V^{2}\right\rangle^{k}\left[\sum_{m} X_{m}(\varepsilon)\right]^{k} \tag{4.19}
\end{equation*}
$$

One sees that the generic term (4.19) for the amplitudes in the two-band problem has exactly the same structure as the term (4.1) with the factor $\left\langle V^{2}\right\rangle X_{m}(\varepsilon)$ playing the role of the factor $V_{0}^{n} V_{n}^{0} /\left(\varepsilon-\Delta_{n}\right)$ and $\varepsilon-\Delta_{n} \rightarrow \varepsilon$. The sum $X_{n}(\varepsilon)$ of the geometric series (4.19) can therefore be given in terms of similar sums $X_{m}(\varepsilon)$, that is

$$
\begin{align*}
X_{n}(\varepsilon) & =\sum_{k} \frac{1}{\left(\varepsilon-\Delta_{n}\right)^{k+1}}\left\langle V^{2}\right\rangle^{k}\left[\sum_{m} X_{m}(\varepsilon)\right]^{k} \\
& =\frac{1}{\varepsilon-\Delta_{n}-\left\langle V^{2}\right\rangle \sum_{m} X_{m}(\varepsilon)} \tag{4.20}
\end{align*}
$$

By analogy

$$
\begin{equation*}
X_{m}(\varepsilon)=\frac{1}{\varepsilon-\Delta_{m}-\left\langle V^{2}\right\rangle \sum_{n} X_{n}(\varepsilon)} \tag{4.21}
\end{equation*}
$$

which together with (4.20) yields a set of two nonlinear algebraic equations

$$
\begin{align*}
& Q_{1}(\varepsilon)=\sum_{n} \frac{1}{\varepsilon-\Delta_{n}-\left\langle V^{2}\right\rangle Q_{2}(\varepsilon)} \\
& Q_{2}(\varepsilon)=\sum_{m} \frac{1}{\varepsilon-\Delta_{m}-\left\langle V^{2}\right\rangle Q_{1}(\varepsilon)} \tag{4.22}
\end{align*}
$$

for two sums $Q_{1}(\varepsilon)=\sum_{n} X_{n}(\varepsilon)$ and $Q_{2}(\varepsilon)=\sum_{m} X_{m}(\varepsilon)$.
Solutions of (4.22) depend on the structure of bands 1 and 2, and several different examples will be considered later in this chapter. A general conclusion emerging from the consideration of the two-band system should already be mentioned: as the result of the perturbation of a two-band system by a matrix of random couplings, the resonant denominators of the fractions $1 /\left(\varepsilon-\Delta_{n ; m}\right)$ that allow for the poles at the positions of energy eigenvalues $\Delta_{n ; m}$ adopt the form $1 /\left(\varepsilon-\Delta_{n ; m}-\left\langle V^{2}\right\rangle Q_{2 ; 1}(\varepsilon)\right)$, which means that all of the energy eigenvalues of the same band acquire an imaginary part $\left\langle V^{2}\right\rangle Q_{1}(\varepsilon)$ or $\left\langle V^{2}\right\rangle Q_{1}(\varepsilon)$, depending on the band number, that are solutions of the set (4.22) of two nonlinear equations. One can refer to this change as renormalization of the energies by a random perturbation that in the regime $t \ll g$ results in the irreversible decay of the states.

Note that by analogy we can obtain the equations

$$
\begin{gather*}
X_{n}(\xi)=\frac{1}{\xi-\Delta_{n}-\left\langle V^{2}\right\rangle Q_{2}(\xi)} ; \quad X_{m}(\xi)=\frac{1}{\xi-\Delta_{m}-\left\langle V^{2}\right\rangle Q_{1}(\xi)} \\
Q_{1}(\xi)=\sum_{n} \frac{1}{\xi-\Delta_{n}-\left\langle V^{2}\right\rangle Q_{2}(\xi)} ; \quad Q_{2}(\xi)=\sum_{m} \frac{1}{\xi-\Delta_{m}-\left\langle V^{2}\right\rangle Q_{1}(\xi)} \tag{4.23}
\end{gather*}
$$

for the complex conjugate amplitudes that coincide with (4.20)-(4.22). However choosing the solution of these equations we should remember that the integration contour of the inverse Fourier transformation over $\xi$ goes not above the real axis, as in Fig. 3.2, but below it, and therefore the solutions for $X_{n ; m}(\xi)$ and for $Q_{1 ; 2}(\xi)$ are the complex conjugates of those for $\varepsilon$.

## Expressions for Averaged Populations

Let us now turn to the perturbation series for the populations. A diagram representing a non-vanishing term of the perturbation series for the populations is shown in Fig. 4.9. Comparing this diagram with the diagrams for a level-band system of Fig. 4.3 one sees a new important feature - the nonvanishing contribution may come from an undulating and dashed lines connecting the same levels $n$ and $m$. Indeed such a diagram contains the product


Fig. 4.9.
$V_{n}^{m}\left(V_{n}^{m}\right)^{*}=V_{n}^{m} V_{m}^{n}$ since the series for the complex conjugate amplitudes have complex conjugate matrix elements, and hence the parallel transition for two different complex conjugated amplitudes yields exactly the same mean value $\left\langle V^{2}\right\rangle$ as two opposite transitions for each of them. In Fig. 4.10 we


Fig. 4.10.
show a higher-order diagram that demonstrates the typical structure of nonvanishing terms - there are two trees "growing" from each level that correspond to the amplitude and conjugate amplitude series, and these trees are connected by coupled lines of different amplitudes.

Employing the notation of (4.20)-(4.23) for the renormalized energies allows one to draw the diagram Fig. 4.10 summed over all possible trees, which takes the form shown in Fig. 4.11


Fig. 4.11.

We can now draw the diagram series for the populations shown in Fig. 4.12.


Fig. 4.12.

All even terms correspond to band 1, and the odd terms correspond to band 2 . One can sum up these perturbation series directly, since they constitute a geometric series with the common factor

$$
\begin{equation*}
\Xi(\varepsilon, \xi)=\sum_{n^{\prime}, m^{\prime}} X_{n^{\prime}}(\varepsilon) X_{n^{\prime}}(\xi)\left\langle V^{2}\right\rangle X_{m}^{\prime}(\varepsilon) X_{m^{\prime}}(\xi)\left\langle V^{2}\right\rangle \tag{4.24}
\end{equation*}
$$

shown in Fig. 4.13.

$$
\Xi(\varepsilon, \xi)=\sum_{n^{\prime}, m^{\prime}} \begin{array}{ccc}
X_{n}^{\prime}(\varepsilon) & X_{m^{\prime}}(\varepsilon) \\
X_{n}(\xi) & X_{m^{\prime}}(\xi)
\end{array}
$$

Fig. 4.13.

The summation yields, for the population $\rho_{m}(\varepsilon, \xi)$ of the states $|m\rangle$,

$$
\begin{align*}
\rho_{m}(\varepsilon, \xi)= & \frac{1}{\varepsilon-\Delta_{m}-V^{2} Q_{1}(\varepsilon)} \frac{1}{\xi-\Delta_{m}-V^{2} Q_{1}(\xi)} \\
& \frac{V^{2}}{1-\Xi(\varepsilon, \xi)} \frac{1}{\xi-V^{2} Q_{2}(\xi)} \frac{1}{\varepsilon-V^{2} Q_{2}(\varepsilon)} \tag{4.25}
\end{align*}
$$

where we write $V^{2}$ instead of $\left\langle V^{2}\right\rangle$ to shorten the notation. We also note an evident albeit remarkable consequence of such a structure of the series

$$
\begin{align*}
\rho_{m}(\varepsilon, \xi)= & \frac{1}{\varepsilon-\Delta_{m}-V^{2} Q_{1}(\varepsilon)} \frac{1}{\xi-\Delta_{m}-V^{2} Q_{1}(\xi)} V^{2} \sum_{n} \rho_{n}(\varepsilon, \xi) \\
\rho_{n}(\varepsilon, \xi)= & \frac{1}{\varepsilon-\Delta_{n}-V^{2} Q_{2}(\varepsilon)} \frac{1}{\xi-\Delta_{n}-V^{2} Q_{2}(\xi)} V^{2} \sum_{m} \rho_{m}(\varepsilon, \xi) \\
& +\delta_{0}^{n} \frac{1}{\xi-V^{2} Q_{2}(\xi)} \frac{1}{\varepsilon-V^{2} Q_{2}(\varepsilon)} \tag{4.26}
\end{align*}
$$

which is depicted in Fig. 4.14 and emerges clearly from the series shown in


Fig. 4.14.

Fig. 4.12. It has an explicit physical meaning - the population of any state of a band forms as a result of transitions from all populated states of the other band.

Equations (4.26) allow one to find the total populations $\rho_{1}(\varepsilon, \xi)=$ $\sum_{n} \rho_{n}(\varepsilon, \xi)$ of the first and $\rho_{2}(\varepsilon, \xi)=\sum_{m} \rho_{m}(\varepsilon, \xi)$ of the second bands after taking the sums over $n$ for the first equation, and over $m$ for the second one. With the allowance of (4.20)-(4.23) this yields the equations

$$
\begin{align*}
& \rho_{1}(\varepsilon, \xi)=\left(V^{2} \sum_{n} X_{n}(\varepsilon) X_{n}(\xi)\right) \rho_{2}(\varepsilon, \xi)+X_{0}(\varepsilon) X_{0}(\xi) \\
& \rho_{2}(\varepsilon, \xi)=\left(V^{2} \sum_{m} X_{m}(\varepsilon) X_{m}(\xi)\right) \rho_{1}(\varepsilon, \xi) \tag{4.27}
\end{align*}
$$

that have the solutions

$$
\begin{align*}
& \rho_{1}(\varepsilon, \xi)=\frac{X_{0}(\varepsilon) X_{0}(\xi)}{1-\left(V^{2} \sum_{n} X_{n}(\varepsilon) X_{n}(\xi)\right)\left(V^{2} \sum_{m} X_{m}(\varepsilon) X_{m}(\xi)\right)} \\
& \rho_{2}(\varepsilon, \xi)=\frac{X_{0}(\varepsilon) X_{0}(\xi)\left(V^{2} \sum_{m} X_{m}(\varepsilon) X_{m}(\xi)\right)}{1-\left(V^{2} \sum_{n} X_{n}(\varepsilon) X_{n}(\xi)\right)\left(V^{2} \sum_{m} X_{m}(\varepsilon) X_{m}(\xi)\right)}, \tag{4.28}
\end{align*}
$$

and for the population difference $\Delta \rho(\varepsilon, \xi)=\rho_{1}(\varepsilon, \xi)-\rho_{2}(\varepsilon, \xi)$ one immediately obtains

$$
\begin{equation*}
\Delta \rho(\varepsilon, \xi)=\frac{\left(1-V^{2} \sum_{m} X_{m}(\varepsilon) X_{m}(\xi)\right) X_{0}(\varepsilon) X_{0}(\xi)}{1-\left(V^{2} \sum_{n} X_{n}(\varepsilon) X_{n}(\xi)\right)\left(V^{2} \sum_{m} X_{m}(\varepsilon) X_{m}(\xi)\right)} \tag{4.29}
\end{equation*}
$$

### 4.2 Non-Degenerate Bands

We are now in the position to consider a particular realization of the two-band system. We start with the simplest example of two statistically equivalent bands, which are the bands with the same mean density of states, shown in Fig. 4.15, and assume that at $t=0$ only the state $n=0$ is populated. We will postpone the important generalization to the case of two different bands until the end of this section, since such a system behaves in qualitatively the same way, although it requires more awkward calculations.

### 4.2.1 General Remarks and the Main Questions

As earlier, we use the energy of the state $|n=0\rangle$ as a reference point for the energy scale, that is we put $\Delta_{n=0}=0$. Note that in order to find a solu-


Fig. 4.15. Two-band system with statistically equivalent bands of identical mean density of states $g_{1}(\Delta)=g_{2}(\Delta)=$ const.
tion for a different initial condition, the superposition principle for the initial populations can be employed due to the fact that the ensemble averaging procedure wipes out all of the interference effects related to the initial conditions. This circumstance is related directly to the tree-like topology of the terms sustaining the ensemble average, which eliminates completely the diagrams with dashed and undulating lines starting at different levels, namely the diagrams responsible for the interference.

We address here two main questions: (i) what are the total populations of the bands as a function of time, and (ii) in which way are these populations distributed among the levels of the bands? The first question can also be posed differently. One can establish a law, which would govern the total population dynamics, and which would have a simple mathematical form.

The second question may have two qualitatively different answers. Indeed, considering the level-band system in the previous chapter, we have seen that the distribution of the populations among the levels of the band follows the

Lorentzian profile (3.35) shown in Fig. 3.5, which has a width given by the level decay rate $W$. For that case the population of a band state $|n\rangle$ has only one possibility to change, namely to return back to the initial level $|0\rangle$, whereas in the two-band system it has other possibilities - the population arriving to a state $|m\rangle$ of the second band from a state $|n\rangle$ of the first band may, in principle return back to any other state $\left|n^{\prime}\right\rangle$ of the band 1 . What happens to the width of the population distribution in this case? Does it become broader by a value $W$ after each subsequent transition from band to band as should be the case if all the phase information is washed away by the ensemble average, or in spite of the ensemble average does the system keep its phase information in such a way that the population distribution over the band remains localized in a finite strip?

### 4.2.2 Renormalized Energies and the Population Distribution

We start by considering the population distribution over the bands and first solve the equations (4.20)-(4.22) for the case $t \ll g$, which allows one to replace the summation by integration. We also note that for the statistically equivalent bands $Q_{1}(\varepsilon)=Q_{2}(\varepsilon)=Q(\varepsilon)$ and $Q_{1}(\xi)=Q_{2}(\xi)=Q(\xi)$, and therefore the equations for $Q(\varepsilon)$ and $Q(\xi)$ take the form

$$
\begin{align*}
& Q(\varepsilon)=\int_{-\Gamma}^{\Gamma} \frac{g d \Delta}{\varepsilon-\Delta-\left\langle V^{2}\right\rangle Q(\varepsilon)} \\
& Q(\xi)=\int_{-\Gamma}^{\Gamma} \frac{g d \Delta}{\xi-\Delta-\left\langle V^{2}\right\rangle Q(\xi)} \tag{4.30}
\end{align*}
$$

where $\Gamma \rightarrow \infty$ is the band width, and $g$ is the state density. We perform the integration and arrive at

$$
\begin{align*}
& Q(\varepsilon)=-g \ln \frac{\Gamma-\left\langle V^{2}\right\rangle Q(\varepsilon)+\varepsilon}{-\Gamma-\left\langle V^{2}\right\rangle Q(\varepsilon)+\varepsilon} \\
& Q(\xi)=-g \ln \frac{\Gamma-\left\langle V^{2}\right\rangle Q(\xi)+\xi}{-\Gamma-\left\langle V^{2}\right\rangle Q(\xi)+\xi} \tag{4.31}
\end{align*}
$$

In the limit $\Gamma \rightarrow \infty$ the numerators of the fractions in the arguments of logarithms are large and positive, and the denominators are large and negative, thus the fraction itself tends to -1 . But $\ln (-1)$ may take both values $\pm i \pi$, depending on the direction of the circumvention of the branching point $x=0$ of $\ln (x)$, which depends in turn on the sign of the imaginary part of the fraction, that is on the locations of zeros of the denominators of $X(\varepsilon)$ and $X(\xi)$ on the complex plane. We remember that $\varepsilon$ has a positive imaginary part, and $\xi$ has a negative one, and therefore the poles of $X(\varepsilon)$ as well as the
zeros of $\Gamma-\left\langle V^{2}\right\rangle Q(\varepsilon)+\varepsilon$ are in the negative imaginary part of the complex plane, whereas for $\xi$ the opposite is the case. This yields in the limit $\Gamma \rightarrow \infty$

$$
\begin{align*}
& Q(\varepsilon)=-i \pi g \\
& Q(\xi)=i \pi g \tag{4.32}
\end{align*}
$$

and hence, according to (4.20)-(4.23)

$$
\begin{align*}
& X_{n}(\varepsilon)=X_{m}(\varepsilon)=X(\varepsilon)=\frac{1}{\varepsilon-\Delta_{n ; m}+i \pi\left\langle V^{2}\right\rangle g} \\
& X_{n}(\xi)=X_{m}(\xi)=X(\xi)=\frac{1}{\xi-\Delta_{n ; m}-i \pi\left\langle V^{2}\right\rangle g} \tag{4.33}
\end{align*}
$$

One sees, that the energy renormalization results in simply adding an imaginary contribution, with the sign changing as a function of the sign of the imaginary part of the variable $\varepsilon$ or $\xi$.

We can now find the sum

$$
\begin{equation*}
\sum_{n} X_{n}(\varepsilon) X_{n}(\xi)=\int_{-\Gamma}^{\Gamma} \frac{g d \Delta}{\left(\varepsilon-\Delta+i \pi\left\langle V^{2}\right\rangle g\right)\left(\xi-\Delta-i \pi\left\langle V^{2}\right\rangle g\right)} \tag{4.34}
\end{equation*}
$$

which enters (4.28) and 4.29). For $\Gamma \rightarrow \infty$ we calculate the residual at the upper part of the complex plane and arrive at

$$
\begin{equation*}
\sum_{n} X_{n}(\varepsilon) X_{n}(\xi)=\frac{-2 \pi i g}{\xi-\varepsilon-2 i \pi\left\langle V^{2}\right\rangle g} \tag{4.35}
\end{equation*}
$$

which also gives the identical value of the sum for the other, statistically equivalent band $\sum_{n} X_{n}(\varepsilon) X_{n}(\xi)=\sum_{m} X_{m}(\varepsilon) X_{m}(\xi)$. After substitution of these sums and (4.32) to (4.25) with the allowance of (4.24) one obtains

$$
\begin{align*}
\rho_{m}(\varepsilon, \xi)= & \frac{1}{\varepsilon-\Delta_{m}+i \pi g V^{2}} \frac{1}{\xi-\Delta_{m}-i \pi g V^{2}} \\
& \frac{V^{2}}{1+\left(\frac{2 \pi g V^{2}}{\xi-\varepsilon-2 i \pi V^{2} g}\right)^{2}} \frac{1}{\xi-i \pi g V^{2}} \frac{1}{\varepsilon+i \pi g V^{2}} \tag{4.36}
\end{align*}
$$

where we as earlier write $V^{2}$ instead of $\left\langle V^{2}\right\rangle$ to shorten the notation. After double inverse Fourier transformation for the time dependent population of the state $|m\rangle$ of the second band we therefore obtain

$$
\begin{align*}
\rho_{m}(t)= & \frac{1}{4 \pi^{2}} \int_{-\infty-i \nu}^{\infty-i \nu} d \xi \int_{-\infty+i \nu}^{\infty+i \nu} d \varepsilon \frac{\exp [i(\xi-\varepsilon) t]}{\left(\xi-i \pi g V^{2}\right)\left(\varepsilon+i \pi g V^{2}\right)(\xi-\varepsilon)} \\
& \frac{V^{2}\left(\xi-\varepsilon-2 i \pi V^{2} g\right)^{2}}{\left(\xi-\varepsilon-4 i \pi V^{2} g\right)\left(\varepsilon-\Delta_{m}+i \pi g V^{2}\right)\left(\xi-\Delta_{m}-i \pi g V^{2}\right)} .(4 \tag{4.37}
\end{align*}
$$

By analogy to (3.113) we introduce the variables $\zeta=\xi-\varepsilon, \eta=(\varepsilon+\xi) / 2$, and perform the integration over $d \eta$, which yields the spectral density of the population $\rho(t, \Delta)=\rho_{m}(t) g$ corresponding to the detuning $\Delta_{m}=\Delta$ :

$$
\begin{equation*}
\rho(t, \Delta)=\frac{1}{\pi^{2}} \int_{-\infty+i \nu}^{\infty+i \nu} d \zeta \frac{\exp [i \zeta t]}{\zeta(\zeta-2 i W)} \frac{W(W+i \zeta)}{(\zeta-i W-\Delta)(\zeta-i W+\Delta)} \tag{4.38}
\end{equation*}
$$

where $W=2 \pi V^{2} g$. Performing integration over $d \zeta$ one obtains the answer to question (ii):

$$
\begin{equation*}
\rho(t, \Delta)=\frac{W}{\pi\left(W^{2}+\Delta^{2}\right)}\left(1+e^{-2 W t}-2 e^{-W t} \cos (\Delta t)\right), \tag{4.39}
\end{equation*}
$$

which indicates clearly that a Lorentzian distribution $W / \pi\left(W^{2}+\Delta^{2}\right)$ of width $W$ attains in the regime $t>1 / W$. Note that for the two-band problem (4.39) gives exactly the same distribution over the levels of the second band as the distribution following from (3.33) for the level-band problem, apart from the fact that the distribution (4.39) is twice as broad and four times smaller than it should be, bearing in mind that for two bands the total sate density is twice as big, and the total population is shared between two bands. We also note that for two statistically equivalent bands the same Lorentzian distribution of populations attains in the first band, although with a slightly different time dependence:

$$
\begin{equation*}
\rho(t, \Delta)=\frac{W}{\pi\left(W^{2}+\Delta^{2}\right)}\left(1-e^{-2 W t}-2 W e^{-W t} \frac{\sin (t \Delta)}{\Delta}\right)+e^{-W t} \delta(\Delta) . \tag{4.40}
\end{equation*}
$$

In Fig. 4.16 we depict the distributions (4.39)-(4.40). One sees that both


Fig. 4.16. Population distribution over the first (a) and the second (b) bands. The initial $\delta$-like distribution of the state 0 is not shown.
population distributions which can be broad initially become localized in the $W$-vicinity of the resonance $\Delta=0$ as time passes.

### 4.2.3 Dynamics of the Total Populations of Bands

Let us now return to question (i) and consider the dynamics of the population distribution between two bands. Equations (4.27) for the total populations with allowance for (4.33) and (4.35) read

$$
\begin{align*}
& \rho_{1}(\varepsilon, \xi)=\frac{-2 \pi i g V^{2}}{\xi-\varepsilon-2 i \pi V^{2} g} \rho_{2}(\varepsilon, \xi)+\frac{1}{\xi-i \pi g V^{2}} \frac{1}{\varepsilon+i \pi g V^{2}} \\
& \rho_{2}(\varepsilon, \xi)=\frac{-2 \pi i g V^{2}}{\xi-\varepsilon-2 i \pi V^{2} g} \rho_{1}(\varepsilon, \xi) \tag{4.41}
\end{align*}
$$

As earlier, by analogy with (3.113) we introduce the variables $\zeta=\varepsilon-\xi$, $\eta=(\varepsilon+\xi) / 2$, and perform the integration over $d \eta$, which yields

$$
\begin{align*}
& \rho_{1}(\zeta)=\frac{-2 \pi i g V^{2}}{\zeta-2 i \pi V^{2} g} \rho_{2}(\zeta)+\frac{i}{\zeta-i 2 \pi g V^{2}} \\
& \rho_{2}(\zeta)=\frac{-2 \pi i g V^{2}}{\varsigma-2 i \pi V^{2} g} \rho_{1}(\zeta) \tag{4.42}
\end{align*}
$$

which for $W=2 \pi V^{2} g$ takes the form

$$
\begin{align*}
& -i \varsigma \rho_{1}(\zeta)=W\left[\rho_{1}(\zeta)-\rho_{2}(\zeta)\right]+1 \\
& -i \varsigma \rho_{2}(\zeta)=W\left[\rho_{2}(\zeta)-\rho_{1}(\zeta)\right] \tag{4.43}
\end{align*}
$$

and after the inverse Fourier transformation over $\zeta$ results in the differential equations

$$
\begin{align*}
& \dot{\rho}_{1}(t)=-W\left[\rho_{1}(t)-\rho_{2}(t)\right]+\delta(t) \\
& \dot{\rho}_{2}(t)=-W\left[\rho_{2}(t)-\rho_{1}(t)\right] \tag{4.44}
\end{align*}
$$

for the time-dependent populations $\rho_{1}(t)$ and $\rho_{2}(t)$ with the $\delta$-function corresponding to the initial condition $\rho_{1}(t=0)=1$.

One immediately recognizes in (4.44) the balance equation, or master equation, and readily finds the population difference

$$
\begin{equation*}
\rho_{1}(t)-\rho_{2}(t)=e^{-2 W t} \tag{4.45}
\end{equation*}
$$

whereas for the sum of the populations one obtains

$$
\begin{equation*}
\dot{\rho}_{1}(t)+\dot{\rho}_{2}(t)=\delta(t), \tag{4.46}
\end{equation*}
$$

which means $\rho_{1}(t)+\rho_{2}(t)=1$ for $t>0$.
Let us summarize the physical meaning of the results obtained. The total populations of the first and the second band obey the master equation, such that if at $t=0$ the population was localized at a single level of the first band, the population difference $\rho_{1}(t)-\rho_{2}(t)$ decreases exponentially from 1 to 0 ,
with a rate $W=4 \pi V^{2} g$ given by the product of the mean squared interaction $V^{2}$ and the state density $g$. After a time $t \sim 1 / W$ both bands have a total population close to $1 / 2$. One may think that this irreversible dynamic is a signature of the loss of coherence in the system. However, the distribution of the populations over the levels of the bands does not become broader with the course of time but on the contrary becomes localized towards a Lorentzian profile of width $W$ which implies that the coherence in the system is present at all times and plays the dominating role in spite of the exponential character of the total population dynamics. The exponential decay is therefore just a transient process that takes place during a limited time interval $t \sim g$, as long as the recurrences do not yet affect the population dynamics. We therefore note an evident condition of the applicability of this approach - there must be many band levels in a strip of the width $W$ of the population distribution. We find this number by multiplying width $W=2 \pi V^{2} g$ by the state density $g$, so that the condition reads $2 \pi V^{2} g^{2} \gg 1$. This also implies that the recurrence time $T=g$ must be much longer as compared to the typical decay time $1 / W$, and that the typical interaction $\sqrt{V^{2}}$ exceeds the mean distance $\delta \sim 1 / g$ among the neighboring band levels. All these conditions allow one to employ the continuous-band model of (4.30).

### 4.2.4 Different Bands

Let us now consider populations $\rho_{n}(\varepsilon, \xi)$ and $\rho_{m}(\varepsilon, \xi)$ of two different bands with the state densities $g_{1}$ and $g_{2}$ respectively. Equations (4.23) yield

$$
\begin{align*}
X_{n}(\varepsilon) & =\frac{1}{\varepsilon-\Delta_{n}-\left\langle V^{2}\right\rangle Q_{2}(\varepsilon)} \\
X_{m}(\varepsilon) & =\frac{1}{\varepsilon-\Delta_{m}-\left\langle V^{2}\right\rangle Q_{1}(\varepsilon)} \\
X_{n}(\xi) & =\frac{1}{\xi-\Delta_{n}-\left\langle V^{2}\right\rangle Q_{2}(\xi)} \\
X_{m}(\xi) & =\frac{1}{\xi-\Delta_{m}-\left\langle V^{2}\right\rangle Q_{1}(\xi)} \tag{4.47}
\end{align*}
$$

where by analogy to (4.30) for the quantities $Q$ we have explicit expressions resulting from the set of equations

$$
\begin{aligned}
& Q_{1}(\varepsilon)=\int_{-\Gamma}^{\Gamma} \frac{g_{1} d \Delta}{\varepsilon-\Delta-\left\langle V^{2}\right\rangle Q_{2}(\varepsilon)}=-i \pi g_{1} \\
& Q_{2}(\varepsilon)=\int_{-\Gamma}^{\Gamma} \frac{g_{2} d \Delta}{\varepsilon-\Delta-\left\langle V^{2}\right\rangle Q_{1}(\varepsilon)}=-i \pi g_{2},
\end{aligned}
$$

$$
\begin{align*}
& Q_{1}(\xi)=\int_{-\Gamma}^{\Gamma} \frac{g_{1} d \Delta}{\xi-\Delta-\left\langle V^{2}\right\rangle Q_{2}(\xi)}=i \pi g_{1} \\
& Q_{2}(\xi)=\int_{-\Gamma}^{\Gamma} \frac{g_{2} d \Delta}{\xi-\Delta-\left\langle V^{2}\right\rangle Q_{1}(\xi)}=i \pi g_{2} \tag{4.48}
\end{align*}
$$

and therefore

$$
\begin{align*}
\sum_{n} X_{n}(\varepsilon) X_{n}(\xi) & =\frac{-2 \pi i g_{1}}{\xi-\varepsilon-2 i \pi\left\langle V^{2}\right\rangle g_{2}} \\
\sum_{m} X_{m}(\varepsilon) X_{m}(\xi) & =\frac{-2 \pi i g_{2}}{\xi-\varepsilon-2 i \pi\left\langle V^{2}\right\rangle g_{1}} \tag{4.49}
\end{align*}
$$

Now (4.28), after integration over $d \eta$, reads

$$
\begin{align*}
\rho_{1}(\zeta) & =\frac{\left(\zeta+i W_{1}\right)}{\left(\zeta+i\left(W_{1}+W_{2}\right)\right) \zeta} \\
\rho_{2}(\zeta) & =\frac{i W_{2}}{\left(\zeta+i\left(W_{1}+W_{2}\right)\right) \zeta} \tag{4.50}
\end{align*}
$$

which yields for the time-dependent and density-weighted population difference

$$
\begin{equation*}
\rho_{1}(t) / g_{1}-\rho_{2}(t) / g_{2}=e^{-\left(W_{1}+W_{2}\right) t} / g_{1} \tag{4.51}
\end{equation*}
$$

Here $W_{1,2}=2 \pi\left\langle V^{2}\right\rangle g_{1,2}$. By analogy to (4.39)-(4.40) a Lorentzian distribution of width $\left(W_{1}+W_{2}\right) / 2$ is attained at $t>1 /\left(W_{1}+W_{2}\right)$. We also note that the total populations of the bands satisfy the master equation

$$
\begin{align*}
\dot{\rho}_{1}(t) & =-W_{2} \rho_{1}(t)+W_{1} \rho_{2}(t)+\delta(t) \\
\dot{\rho}_{2}(t) & =-W_{1} \rho_{2}(t)+W_{2} \rho_{1}(t) \tag{4.52}
\end{align*}
$$

### 4.3 Two Degenerate Levels

The coherent character of the population dynamics discussed at the end of Sect. 4.2.3 becomes much more apparent for the case of two degenerate levels, that is for bands with $\Delta_{m}=\Delta_{n}=0$. Considering such a system we first make use of the results of the general analysis of Sect. 4.1, and then we present an alternative, that allows one to obtain a deeper insight into the random matrix model.

### 4.3.1 Degenerate Levels as a Complex System

Equations (4.20)-(4.23), (4.26)-(4.28) remain valid for this system, although the particular algebraic form of (4.20)-(4.23) is now different from (4.32)(4.33). Let us assume, as earlier in (4.30), that the bands are statistically
equivalent both having an identical number $N$ of levels, and write (4.22) for $Q_{1,2}(\varepsilon)=Q(\varepsilon)$ and $Q_{1,2}(\xi)=Q(\xi)$ in the form

$$
\begin{align*}
Q(\varepsilon) & =\frac{N}{\varepsilon-\left\langle V^{2}\right\rangle Q(\varepsilon)} \\
Q(\xi) & =\frac{N}{\xi-\left\langle V^{2}\right\rangle Q(\xi)} \tag{4.53}
\end{align*}
$$

These second-order algebraic equations yield

$$
\begin{align*}
& Q(\varepsilon)=\frac{\varepsilon+\sqrt{\varepsilon^{2}-4 N\left\langle V^{2}\right\rangle}}{2\left\langle V^{2}\right\rangle} \\
& Q(\xi)=\frac{\xi-\sqrt{\xi^{2}-4 N\left\langle V^{2}\right\rangle}}{2\left\langle V^{2}\right\rangle} \tag{4.54}
\end{align*}
$$

where the signs of the roots are chosen such that the signs of the imaginary parts of $Q(\varepsilon)$ and $Q(\xi)$ coincide with those for $\varepsilon^{-1}$ and $\xi^{-1}$ on the integration contours ( $-\infty \pm i \nu ; \infty \pm i \nu$ ) of the inverse Fourier transformation respectively, since otherwise the inverse Fourier transforms of the corresponding probability amplitudes vanish. Equations (4.32) give us a hint for the proper choice.

All states of each of the degenerate bands are apparently statistically equivalent, and therefore we will not consider the population distribution amongst them but dwell on the difference of the total populations of the bands. Substitution of (4.54) into (4.20), (4.23) yields

$$
\begin{align*}
& X_{n}(\varepsilon)=X_{m}(\varepsilon)=X(\varepsilon)=\frac{2}{\varepsilon-\sqrt{\varepsilon^{2}-4 N\left\langle V^{2}\right\rangle}} \\
& X_{n}(\xi)=X_{m}(\xi)=X(\xi)=\frac{2}{\xi+\sqrt{\xi^{2}-4 N\left\langle V^{2}\right\rangle}} \tag{4.55}
\end{align*}
$$

and therefore (4.29) results in

$$
\begin{equation*}
\Delta \rho(\varepsilon, \xi)=\frac{X(\varepsilon) X(\xi)}{1+N\left\langle V^{2}\right\rangle X(\varepsilon) X(\xi)} \tag{4.56}
\end{equation*}
$$

The inverse Fourier transform of (4.56) reads

$$
\begin{equation*}
\Delta \rho(t)=\frac{1}{4 \pi^{2}} \int_{-\infty+i \nu}^{\infty+i \nu} \int_{-\infty-i \nu}^{\infty-i \nu} \frac{\exp [i(\varepsilon-\xi) t] d \varepsilon d \xi}{\left(\varepsilon-\sqrt{\varepsilon^{2}-4 N V^{2}}\right)\left(\xi+\sqrt{\xi^{2}-4 N V^{2}}\right)+N V^{2}}, \tag{4.57}
\end{equation*}
$$

where as earlier we replace $\left\langle V^{2}\right\rangle$ by $V^{2}$.
One can close the integration contours in the complex planes $\varepsilon$ and $\xi$ and transform them such that they become the closed loops $C_{1}$ and $C_{2}$ going
around the branching points $\varepsilon_{1,2}=\xi_{1,2}= \pm \sqrt{4 N V^{2}}$ as shown in Fig. 4.17(a). After the replacements $\varepsilon=\sqrt{4 N V^{2}} \sin \varphi$ and $\xi=\sqrt{4 N V^{2}} \sin \theta$ these contours become intervals as shown in Fig. 4.17(b), whereas the integral (4.57)


Fig. 4.17. Integration contours $C_{1}$ and $C_{2}$ in the complex planes $\varepsilon$ and $\xi$ respectively (a) go around the points $w= \pm \sqrt{4 N V^{2}}$. The replacements $\varepsilon \rightarrow \sqrt{4 N V^{2}} \sin \varphi$ and $\xi \rightarrow \sqrt{4 N V^{2}} \sin \theta$ transform these contours to the straight lines (b) in the complex plane of $\varphi$ and $\theta$
takes the form

$$
\begin{equation*}
\Delta \rho(t)=\int_{C_{1}} \int_{C_{2}} \frac{\exp \left[i \sqrt{4 N V^{2}}(\sin \varphi-\sin \theta) t\right] \cos \varphi \cos \theta d \varphi d \theta}{\pi^{2}\left(e^{i(\varphi-\theta)}+1\right)} \tag{4.58}
\end{equation*}
$$

We first perform the integration over one of the variables, say $\varphi$, by moving the integration contour $C_{1}$ toward $\varphi \rightarrow i \infty$ along the strip ( $-\pi, \pi$ ), by taking the residual at the point $\varphi=\theta+\pi$, and allowing for the fact that as $\varphi \rightarrow i \infty$ the integrand vanishes, we arrive at

$$
\begin{equation*}
\Delta \rho(t)=\frac{2}{\pi} \int_{-\pi}^{\pi} d \theta \exp \left[-i 4 t \sqrt{N V^{2}} \sin \theta\right] \cos ^{2} \theta \tag{4.59}
\end{equation*}
$$

The integral (4.59) can be given in terms of Bessel functions. Indeed

$$
\begin{equation*}
J_{n}(z)=\frac{1}{2 \pi} \int_{-\pi}^{\pi} d \theta \exp [i n \theta+i z \sin \theta] \tag{4.60}
\end{equation*}
$$

and $2 \cos ^{2} \theta=1+\cos 2 \theta$, and therefore

$$
\begin{equation*}
\Delta \rho(t)=J_{0}\left(4 t \sqrt{N V^{2}}\right)+J_{2}\left(4 t \sqrt{N V^{2}}\right)=\frac{J_{1}\left(4 t \sqrt{N V^{2}}\right)}{2 t \sqrt{N V^{2}}} \tag{4.61}
\end{equation*}
$$

In Fig. 4.18 we show this dependence. One can see that the population difference of the first and the second bands is an oscillatory and decaying function of time. This means that the populations of these two degenerate levels become equal, that is they indeed do tend to the limit given by the irreversible master equation, although this process does not have an exponential character, but on the contrary, has strong reminiscences of coherent behavior. Therefore such dynamics are sometimes called coherent damping.


Fig. 4.18. Two degenerate levels (a). The arrows show possible transitions and correspond to different transition matrix elements $V_{n}^{m}$. (b) Population difference of two statistically equivalent degenerate bands as a function of dimensionless time $\tau=$ $t \sqrt{N V^{2}}$. The population dynamics is not exponential, and is influenced strongly by the coherence, although at the limit of long times it results in the equal distribution of the population among the bands, typical of incoherent decay, as one can see from the total populations $\rho_{1}$ and $\rho_{2}$ shown in (c).

### 4.3.2 The Bands as an Ensemble of Two-Level Systems

Let us consider the results obtained from a different point of view. The Hamiltonian $\hat{H}$ of the system of two degenerate levels coupled to each other by a random matrix has a specific block structure:

$$
\hat{H}=\left(\begin{array}{cc}
\hat{0} & \hat{V}  \tag{4.62}\\
\hat{V}^{+} & \hat{0}
\end{array}\right)=\left(\begin{array}{cc}
0 & V_{n}^{m} \\
V_{m}^{n} & 0
\end{array}\right)
$$

with zero diagonal blocks. The zero matrix has an evident property - it remains zero after multiplication by any other matrix. We can therefore construct a unitary matrix

$$
\hat{S}=\left(\begin{array}{cc}
\hat{s} & \hat{0}  \tag{4.63}\\
\hat{0} & \hat{s}^{+}
\end{array}\right)=\left(\begin{array}{cc}
s_{n}^{n^{\prime}} & 0 \\
0 & \left(s^{+}\right)_{m}^{m^{\prime}}
\end{array}\right)
$$

where the unitary submatrix $s_{n}^{n^{\prime}}$ and its conjugate $\left(s^{+}\right)_{m}^{m^{\prime}}$ diagonalize the matrix $V_{n}^{m}$, that is $\sum_{n^{\prime} m^{\prime}} s_{n}^{n^{\prime}} V_{n^{\prime}}^{m^{\prime}}\left(s^{+}\right)_{m^{\prime}}^{m}=V_{n} \delta_{n}^{m}$, where $V_{n}$ are the eigenvalues of $\hat{V}$. This implies that the matrix $\hat{V}$ is Hermitian, which is a priori not necessarily the case although we assume it in this section; and in the general case the resulting matrix is not diagonal but a triangular one. After the transformation $\hat{S} \hat{H} \hat{S}^{+}$the Hamiltonian takes the form

$$
\hat{H}=\left(\begin{array}{cc}
0 & V_{n} \delta_{n}^{m}  \tag{4.64}\\
V_{n} \delta_{n}^{m} & 0
\end{array}\right)
$$

with diagonal matrices in the off-diagonal blocks. This form has a physical meaning: we have represented our system of two degenerate levels coupled by
a random matrix as a set of $N$ independent two-level systems each of which has the coupling matrix element $V_{n}$. The random character of the coupling also suggests that in our new representation all the states corresponding to the level 1 are uniformly populated.

In Chap. 3 we have seen that the population of the upper level of a resonant two-level system with the coupling $V$ is given by (3.20) and has the form $\rho_{u p}(t)=\sin ^{2} V t$, whereas the population of the lower level is $\rho_{l w}(t)=$ $\cos ^{2} V t$, and hence the population difference $\Delta \rho_{n}(t)$ corresponding to the coupling $V_{n}$ reads

$$
\begin{equation*}
\Delta \rho_{n}(t)=\cos ^{2} V_{n} t-\sin ^{2} V_{n} t=\cos 2 V_{n} t \tag{4.65}
\end{equation*}
$$

The ensemble average population is therefore given by (4.65) averaged with the distribution function of the couplings $g\left(V_{n}\right)$, that is

$$
\begin{equation*}
\Delta \rho(t)=\int d V_{n} g\left(V_{n}\right) \cos 2 V_{n} t \tag{4.66}
\end{equation*}
$$

From (4.66) we immediately obtain (4.61) by taking the distribution

$$
\begin{equation*}
g\left(V_{n}\right)=\frac{\sqrt{\left(V_{n}\right)^{2}-4 N\left\langle V^{2}\right\rangle}}{2 \pi N\left\langle V^{2}\right\rangle} \tag{4.67}
\end{equation*}
$$

of the span $\pm 2 \sqrt{N\left\langle V^{2}\right\rangle}$, which is known as the Wigner semicircular distribution (1.2) of the eigenvalues of a random matrix with Gaussian statistics of the matrix elements $g\left(V_{n}^{m}\right)=e^{-\left(V_{n}^{m}\right)^{2} /\left\langle V^{2}\right\rangle}$. Indeed, when we substitute (4.67) into (4.66) and make the replacement $V_{n} \rightarrow 2 \sqrt{N\left\langle V^{2}\right\rangle} \sin \theta$ we arrive at

$$
\begin{equation*}
\Delta \rho(t)=\frac{2}{\pi} \int_{-\pi}^{\pi} d \theta \cos ^{2} \theta \cos \left[4 t \sqrt{N\left\langle V^{2}\right\rangle} \sin \theta\right] \tag{4.68}
\end{equation*}
$$

which coincides with (4.59) if we take parity into account. This coincidence is in fact an independent confirmation of the Wigner distribution given by (4.67). We also note that the imaginary part $\operatorname{Im}\left[\frac{X(\varepsilon)}{\pi}\right]=$ $\frac{2}{\pi} \operatorname{Im}\left[\varepsilon-\sqrt{\varepsilon^{2}-4 N\left\langle V^{2}\right\rangle}\right]=\sqrt{4 N\left\langle V^{2}\right\rangle-\varepsilon^{2}} / 2 \pi N\left\langle V^{2}\right\rangle$ of the average resolvent represents the density of eigenstates at the energy $E=\varepsilon$. This part also yields the distribution (4.67), which is quite natural since the two energy eigenstates of a resonant two-level system (3.14) are at the points $E_{1,2}= \pm V_{n}$ given by the value of the level coupling $V_{n}$.

In Sect. 4.1, when we were considering the role of returns, we mentioned that the approach based on the assumption of the tree-like topology of the diagrams may fail when the number of transitions exceeds the number of levels involved in the process and the self-intersections of the diagrams mentioned on p. 131 start to play a role. This is also the case for the degenerate bands. One estimates the number of transitions during the time $t$ by the product


Fig. 4.19. Spectrum of eigenvalues of a $100 \times 100$ random Hermitian matrix with Gaussian distribution of the matrix elements shown by ticks on the abscissa, and the convolution of this spectrum with a Gaussian of width 0.5 (solid line) that gives the averaged state density. The corresponding Wigner distribution of the eigenvalues of a random matrix is shown by the dashed line. One sees that the eigenvalue spectrum is far from being equidistant.
$\sqrt{\left\langle V^{2}\right\rangle} t$ and the number of levels is apparently $N$. This means that the approach is valid as long as $t<N / \sqrt{\left\langle V^{2}\right\rangle}$. This time exceeds considerably the mean energy eigenstate density $g \sim \sqrt{N /\left\langle V^{2}\right\rangle}$ of the coupled degenerate bands, which can be estimated as the total number of states $2 N$ divided by the width $2 \sqrt{N\left\langle V^{2}\right\rangle}$ of the distribution of the couplings (4.67).

It is expedient to consider this situation in the context of the quantum recurrences and quantum revivals discussed in Sect. 3.3 that are the physical phenomena destroying the validity of the tree-like topology approximation. If the exact energy eigenstate spectrum is equidistant, or close to equidistant, then at $t \sim g$ one observes a strong recurrence of the population to the initial state that cannot be found in the framework of the continuous-band model. But if the eigenstates are randomly placed on the energy scale, there is no recurrence at $t \sim g$, although a partial revival may take place at a longer time. The latter is rather the case for the system of two degenerate levels, where the spectrum of the coupling matrix elements $V_{n}$ has no physical reasons to be equidistant for small values of $V$. We illustrate this in Fig. 4.19 where an example of a random matrix eigenstates is shown along with the Wigner semicircular distribution.

### 4.4 A Band Coupled to a Degenerate Level

We consider now one more extreme, namely a system of two different bands: an $N$-fold degenerate level 1 coupled to a band 2 with a constant density of states $g$. This system has both the features typical of the exponential dynamics of a two-non-degenerate band system and the features typical of the coherent dynamics system of two degenerate levels.

### 4.4.1 Total Population of the Bands

Equations (4.22) for the case of one degenerate and one uniform band reads

$$
\begin{align*}
Q_{1}(\varepsilon) & =\frac{N}{\varepsilon-\left\langle V^{2}\right\rangle Q_{2}(\varepsilon)} \\
Q_{2}(\varepsilon) & =\sum_{m} \frac{1}{\varepsilon-\Delta_{m}-\left\langle V^{2}\right\rangle Q_{1}(\varepsilon)} \tag{4.69}
\end{align*}
$$

and one finds the equations for $Q_{1,2}(\xi)$ analogously. We consider the case $t \ll$ $g$ and by analogy to (4.48) replace the sum by the integral, which immediately yields

$$
\begin{align*}
Q_{1}(\varepsilon) & =\frac{N}{\varepsilon+i \pi\left\langle V^{2}\right\rangle g} \\
Q_{2}(\varepsilon) & =-i \pi g \tag{4.70}
\end{align*}
$$

and

$$
\begin{align*}
Q_{1}(\xi) & =\frac{N}{\xi-i \pi\left\langle V^{2}\right\rangle g} \\
Q_{2}(\xi) & =i \pi g \tag{4.71}
\end{align*}
$$

We also note that

$$
\begin{align*}
X_{n}(\varepsilon) & =\frac{1}{\varepsilon+i \pi\left\langle V^{2}\right\rangle g} \\
X_{m}(\varepsilon) & =\frac{1}{\varepsilon-\Delta_{m}-\left\langle V^{2}\right\rangle N\left(\varepsilon+i \pi\left\langle V^{2}\right\rangle g\right)^{-1}} \\
X_{n}(\xi) & =\frac{1}{\xi-i \pi\left\langle V^{2}\right\rangle g} \\
X_{m}(\xi) & =\frac{1}{\xi-\Delta_{m}-N\left\langle V^{2}\right\rangle\left(\xi-i \pi\left\langle V^{2}\right\rangle g\right)^{-1}} \tag{4.72}
\end{align*}
$$

and hence

$$
\begin{align*}
\sum_{n} X_{n}(\varepsilon) X_{n}(\xi) & =\frac{N}{\left(\varepsilon+i \pi\left\langle V^{2}\right\rangle g\right)\left(\xi-i \pi\left\langle V^{2}\right\rangle g\right)} \\
\sum_{m} X_{m}(\varepsilon) X_{m}(\xi) & =\frac{-2 \pi i g}{\xi-\varepsilon-\left\langle V^{2}\right\rangle\left(Q_{1}(\xi)-Q_{1}(\varepsilon)\right)} \tag{4.73}
\end{align*}
$$

where

$$
\begin{equation*}
Q_{1}(\xi)-Q_{1}(\varepsilon)=\frac{N\left(\varepsilon-\xi+2 i \pi\left\langle V^{2}\right\rangle g\right)}{\left(\varepsilon+i \pi\left\langle V^{2}\right\rangle g\right)\left(\xi-i \pi\left\langle V^{2}\right\rangle g\right)} \tag{4.74}
\end{equation*}
$$

We substitute (4.73), (4.74) into (4.29) and after straightforward algebraic transformations arrive at

$$
\begin{align*}
\Delta \rho(\varepsilon, \xi)= & \left(\frac{1}{\varepsilon+i \pi\left\langle V^{2}\right\rangle g}-\frac{1}{\xi-i \pi\left\langle V^{2}\right\rangle g}\right) \frac{1}{\xi-\varepsilon}  \tag{4.75}\\
& +\frac{4 \pi i g\left\langle V^{2}\right\rangle}{\left(\varepsilon+i \pi\left\langle V^{2}\right\rangle g\right)\left(\xi-i \pi\left\langle V^{2}\right\rangle g\right)+\left\langle V^{2}\right\rangle N} \frac{1}{\xi-\varepsilon}
\end{align*}
$$

We now note that $i(\xi-\varepsilon) \Delta \rho(\varepsilon, \xi)$ is the Fourier transform of the time derivative of the population difference, which we denote $\Delta \rho^{\prime}(\varepsilon, \xi)$, and obtain

$$
\begin{align*}
\Delta \rho^{\prime}(\varepsilon, \xi)= & \frac{i}{\left(\varepsilon+i \pi\left\langle V^{2}\right\rangle g\right)}-\frac{i}{\left(\xi-i \pi\left\langle V^{2}\right\rangle g\right)} \\
& -\frac{4 \pi g\left\langle V^{2}\right\rangle}{\left(\varepsilon+i \pi\left\langle V^{2}\right\rangle g\right)\left(\xi-i \pi\left\langle V^{2}\right\rangle g\right)+\left\langle V^{2}\right\rangle N} \tag{4.76}
\end{align*}
$$

The first two terms depend only on one of the two variables $\varepsilon$ and $\xi$, and therefore for $t>0$ they vanish after the inverse Fourier transformation over both $\varepsilon$ and $\xi$. Hence we have

$$
\begin{equation*}
\frac{d}{d t} \Delta \rho(t)=\frac{-1}{4 \pi^{2}} \int_{-\infty+i \nu}^{\infty+i \nu} \int_{-\infty-i \nu}^{\infty-i \nu} \frac{4 \pi i g\left\langle V^{2}\right\rangle \exp [i(\xi-\varepsilon) t] d \varepsilon d \xi}{\left(\varepsilon+i \pi\left\langle V^{2}\right\rangle g\right)\left(\xi-i \pi\left\langle V^{2}\right\rangle g\right)+\left\langle V^{2}\right\rangle N} \tag{4.77}
\end{equation*}
$$

We now introduce the variables $\zeta=\xi-\varepsilon, \eta=(\varepsilon+\xi) / 2$, take into account the relation $\left(\varepsilon+i \pi\left\langle V^{2}\right\rangle g\right)\left(\xi-i \pi\left\langle V^{2}\right\rangle g\right)=\eta^{2}-\left(\zeta / 2-i \pi\left\langle V^{2}\right\rangle g\right)^{2}$, and perform the integration over $d \eta$, which yields

$$
\begin{equation*}
\frac{d}{d t} \Delta \rho(t)=g\left\langle V^{2}\right\rangle \int_{-\infty+i \nu}^{\infty+i \nu} \frac{\exp [i \zeta t] d \zeta}{\sqrt{\left(\zeta / 2-i \pi\left\langle V^{2}\right\rangle g\right)^{2}-\left\langle V^{2}\right\rangle N}} \tag{4.78}
\end{equation*}
$$

One can close the integration path, such that it becomes a closed contour going around the interval $\left(i 2 \pi\left\langle V^{2}\right\rangle g-2 \sqrt{\left\langle V^{2}\right\rangle N}, i 2 \pi\left\langle V^{2}\right\rangle g+2 \sqrt{\left\langle V^{2}\right\rangle N}\right)$ connecting two branching points in the complex plane $\zeta$, and obtain after the replacement $\zeta \rightarrow i 2 \pi\left\langle V^{2}\right\rangle g+2 \sqrt{\left\langle V^{2}\right\rangle N} \cos \varphi$

$$
\begin{align*}
\frac{d}{d t} \Delta \rho(t) & =-4 g\left\langle V^{2}\right\rangle \exp \left[-2 \pi g\left\langle V^{2}\right\rangle t\right] \int_{-\pi}^{\pi} \exp \left[i t 2 \sqrt{\left\langle V^{2}\right\rangle N} \cos \varphi\right] d \varphi \\
& =-4 \pi g\left\langle V^{2}\right\rangle \exp \left[-2 \pi g\left\langle V^{2}\right\rangle t\right] J_{0}\left(t \sqrt{4\left\langle V^{2}\right\rangle N}\right) \tag{4.79}
\end{align*}
$$

that is

$$
\begin{equation*}
\Delta \rho(t)=1-4 \pi g\left\langle V^{2}\right\rangle \int_{0}^{t} \exp \left[-2 \pi\left\langle V^{2}\right\rangle g \tau\right] J_{0}\left(\tau \sqrt{4\left\langle V^{2}\right\rangle N}\right) d \tau \tag{4.80}
\end{equation*}
$$

where we have taken into account the initial condition $\Delta \rho(t=0)=1$.
One sees clearly the oscillatory character of the time derivative of the population difference coming from the Bessel function in (4.79), and some reminiscences of these oscillations are seen in the dependence of the population difference $\Delta \rho(t)$ of (4.80) itself, which is shown in Fig. 4.20. This behav-


Fig. 4.20. The population difference $\rho_{1}(t)-\rho_{2}(t)$ for a system of an $N$-fold degenerate level coupled to a band with a state density $g$ by a random matrix with the mean squared matrix element $\left\langle V^{2}\right\rangle$ as a function of the scaled time $T=t \sqrt{4\left\langle V^{2}\right\rangle N}$ is shown (a) for several values of the dimensionless parameter $w=\pi g \sqrt{\left\langle V^{2}\right\rangle / N}$. The asymptotic values of the populations of the degenerate level $\rho_{1}(\infty)$ and of the band $\rho_{2}(\infty)$ depend (b) on the parameter $w$.
ior is typical of the coherent damping in the system of two non-degenerate bands. At the same time, in (4.79) one also sees the exponential factor $\exp \left[-2 \pi g\left\langle V^{2}\right\rangle t\right]$ typical of the level-band problem and for the two nondegenerate bands problem. And we see here once more the interplay of the coherent and relaxation processes in the population dynamics of the complex multilevel systems, although we have to remember that the exponential relaxation occurs only until the quantum revivals and recurrences start to play a role at $t>g$.

Let us concentrate on the populations at times $t \gg 1 / \sqrt{4 N\left\langle V^{2}\right\rangle}$ and $t \gg 1 / 2 \pi g\left\langle V^{2}\right\rangle$, assuming that $t \ll g$. At this extreme we can extend the integration limit in (4.80) to infinity and obtain

$$
\begin{equation*}
\Delta \rho(\infty)=1-\frac{4 \pi g\left\langle V^{2}\right\rangle}{\sqrt{\left(2 \pi\left\langle V^{2}\right\rangle g\right)^{2}+4\left\langle V^{2}\right\rangle N}} \tag{4.81}
\end{equation*}
$$

where we have taken into account that $\int_{0}^{\infty} \exp [-\alpha \tau] J_{0}(\tau \beta) d \tau=\left(\alpha^{2}+\right.$ $\left.\beta^{2}\right)^{-1 / 2}$. If we make use of the relation $\rho_{1}(\infty)+\rho_{2}(\infty)=1$ we get immediately
the explicit asymptotic populations

$$
\begin{align*}
& \rho_{1}(\infty)=1-\frac{2 \pi g\left\langle V^{2}\right\rangle}{\sqrt{\left(2 \pi\left\langle V^{2}\right\rangle g\right)^{2}+4\left\langle V^{2}\right\rangle N}}, \\
& \rho_{2}(\infty)=\frac{2 \pi g\left\langle V^{2}\right\rangle}{\sqrt{\left(2 \pi\left\langle V^{2}\right\rangle g\right)^{2}+4\left\langle V^{2}\right\rangle N}} \tag{4.82}
\end{align*}
$$

shown in Fig. 4.20(b). The total populations of the degenerate level and the band become equal for a mean squared coupling $\left\langle V^{2}\right\rangle=V_{s t}^{2}=N / \pi^{2} g^{2}$ for which the number of the populated energy eigenstates of the band equals the degeneracy of the level.

### 4.4.2 Population Distribution over the Band

Let us concentrate on the population distribution among the states of the band. After straightforward, albeit cumbersome transformations, (4.26) with allowance of (4.28) and (4.71)-(4.74) yields

$$
\begin{align*}
\rho_{m}(\varepsilon, \xi)= & \frac{1}{\left(\varepsilon-\Delta_{m}\right)(\varepsilon+i W)-\frac{W^{2}}{w^{2}}} \frac{1}{\left(\xi-\Delta_{m}\right)(\xi-i W)-\frac{W^{2}}{w^{2}}} \\
& {\left[1-\frac{2 i W^{3} / w^{2}}{\left[(\varepsilon+i W)(\xi-i W)+\frac{W^{2}}{w^{2}}\right](\xi-\varepsilon)}\right] \frac{W^{2}}{N w^{2}} } \tag{4.83}
\end{align*}
$$

where we denote $W=\pi\left\langle V^{2}\right\rangle g$ and $w=W / \sqrt{\left\langle V^{2}\right\rangle N}$. Further analysis of this expression should rely on the search for the positions of poles in the complex planes of the variables $\varepsilon$ and $\xi$ followed by the inverse Fourier transformation with respect to these variables, which also require quite cumbersome calculations. However, these calculations become much shorter if we concentrate only on the distribution which is attained in the regime $t \rightarrow \infty$, when the only non-vanishing contribution is given by the second term in the parentheses in the second line of (4.83) containing the factor $(\xi-\varepsilon)$ in the denominator. This contribution gives a time independent term for the population distribution $\rho_{m}(t=\infty)=\rho_{m}$ which reads

$$
\begin{align*}
\rho_{m}= & \frac{1}{\pi} \int_{-\infty}^{\infty} \frac{1}{\left[\left(\varepsilon-\Delta_{m}\right)(\varepsilon+i W)-\frac{W^{2}}{w^{2}}\right]} \frac{1}{\left[\left(\varepsilon-\Delta_{m}\right)(\varepsilon-i W)-\frac{W^{2}}{w^{2}}\right]} \\
& \frac{W^{5} / w^{4} N}{\left[(\varepsilon+i W)(\varepsilon-i W)+\frac{W^{2}}{w^{2}}\right]} d \varepsilon . \tag{4.84}
\end{align*}
$$

The integrand (4.84) has six poles at the points of the complex plane $\varepsilon$ located symmetrically with respect to the real axis. By taking residuals at
three points of the upper part of the complex plane and factoring out $W^{2}$ with allowance for the relation $N w=\pi g W / w$ one obtains

$$
\begin{equation*}
\rho_{m} g=\frac{1}{\pi W} \frac{w}{\sqrt{w^{2}+1}} \frac{2 \sqrt{1+1 / w^{2}}-1}{\left(2 \sqrt{1+1 / w^{2}}-1\right)^{2}+\left(\frac{\Delta_{m}}{W}\right)^{2}} . \tag{4.85}
\end{equation*}
$$

In (4.85) one recognizes the total population $\rho_{2}(\infty)=w / \sqrt{w^{2}+1}$ of the non-degenerate band (4.82) multiplied by a Lorentzian distribution of width $\Gamma=2 \sqrt{1+1 / w^{2}}-1$ over scaled detunings $\delta=\Delta_{m} / W$, which reads

$$
\begin{align*}
\rho_{m}= & \frac{\left\langle V^{2}\right\rangle}{\sqrt{\left(\pi\left\langle V^{2}\right\rangle g\right)^{2}+\left\langle V^{2}\right\rangle N}} \\
& \frac{2 \sqrt{\left(\pi\left\langle V^{2}\right\rangle g\right)^{2}+\left\langle V^{2}\right\rangle N}-\pi\left\langle V^{2}\right\rangle g}{\left(2 \sqrt{\left(\pi\left\langle V^{2}\right\rangle g\right)^{2}+\left\langle V^{2}\right\rangle N}-\pi\left\langle V^{2}\right\rangle g\right)^{2}+\Delta_{m}^{2}}, \tag{4.86}
\end{align*}
$$

in the initial notation whereas the width is $\Gamma=2 \sqrt{\left(\pi\left\langle V^{2}\right\rangle g\right)^{2}+\left\langle V^{2}\right\rangle N}-$ $\pi\left\langle V^{2}\right\rangle g$. This dependence (4.85) is shown in Fig. 4.21(a), and we replace hereafter in this subsection $\Delta_{m}$ by $\Delta$. One sees that for large coupling,


Fig. 4.21. The asymptotic population distribution (a) over the levels of the band as a function of the scaled detuning $\Delta / \pi g\left\langle V^{2}\right\rangle$ and the dimensionless parameter $w=$ $\pi g \sqrt{\left\langle V^{2}\right\rangle / N}$. The asymptotic distribution is attained as a result of a dynamical process that yields a distribution oscillating in time shown in (b) as a function of the detuning $\Delta$ for $w=1 / \sqrt{3}$.
$\left\langle V^{2}\right\rangle \rightarrow \infty$, when $\sqrt{\left\langle V^{2}\right\rangle N}$ becomes small as compared to $W=\pi g\left\langle V^{2}\right\rangle$ and as a consequence the number $\pi g^{2}\left\langle V^{2}\right\rangle$ of the band levels participating in the population dynamics considerably exceeds the degeneracy $N$ of the
level, and the width $\Gamma$ of the Lorentzian distribution equals $W=\pi\left\langle V^{2}\right\rangle g$ as is the case for the level-band problem (3.35) discussed in Chap. 3. This is quite natural, since for each state of the degenerate level, there are enough energy eigenstates of the band where its population can go without interfering with the decay of other states of the level. In the opposite limit $\sqrt{\left\langle V^{2}\right\rangle N}$ $\gg \pi g\left\langle V^{2}\right\rangle$ the distribution width is $\Gamma \simeq 2 \sqrt{\left\langle V^{2}\right\rangle N}$, and the number of band levels in resonance amounts to $\pi \Gamma g \simeq 2 \pi \sqrt{\left\langle V^{2}\right\rangle N g^{2}}$. Note that if we assume that the total populations $\rho_{2}$ of the band and $\rho_{1}$ of the degenerate level are proportional to the corresponding numbers of states in resonance, we arrive at $\rho_{2} / \rho_{1}=2 \pi \sqrt{\left\langle V^{2}\right\rangle N g^{2}} / N=2 \pi\left\langle V^{2}\right\rangle g / \sqrt{\left\langle V^{2}\right\rangle N}$, which is twice as large as the value suggested in the same extreme by the exact solution (4.82). This means that the equal distribution principle does not always govern the stationary population distribution among the states of a complex multilevel system.

In order to trace the dynamics of the population distribution over the band of levels one needs to perform the inverse Fourier transformation of (4.83), which requires straightforward although cumbersome calculations. We do not present here all the details, but just show the crucial aspects of the main transformation, that could serve as guiding lines for such calculations. The unity in the square brackets yields just an absolute value of the inverse Fourier transform of

$$
\begin{equation*}
\psi_{m}(\varepsilon)=\frac{1}{(\varepsilon-\Delta)(\varepsilon+i W)-\frac{W^{2}}{w^{2}}} \sqrt{\frac{W^{2}}{N w^{2}}} \tag{4.87}
\end{equation*}
$$

which is not the total ensemble average population, but "a naive population" $\tilde{\rho}_{m}(t, \Delta)$ given by the absolute value squared of the ensemble averaged wavefunction. The right-hand side of (4.87) has poles at the points

$$
\begin{equation*}
\varepsilon_{1,2}=\frac{1}{2}(\Delta-i W) \pm \sqrt{\frac{1}{2}(\Delta-i W)^{2}+\frac{W^{2}}{w^{2}}} \tag{4.88}
\end{equation*}
$$

and the integration performed by calculating the residuals at these points gives the time-dependent average amplitude, which yields the corresponding "probability"

$$
\begin{equation*}
\widetilde{\rho}_{m}(t, \Delta)=\frac{W^{2}}{N w^{2}} \frac{2\left|\sin \left(\frac{t}{2} \sqrt{4 \frac{W^{2}}{w^{2}}+(\Delta-i W)^{2}}\right)\right|^{2} e^{-t W}}{\sqrt{\Delta^{4}+\left(W^{2}-4 \frac{W^{2}}{w^{2}}\right)^{2}+2 \Delta^{2}\left(4 \frac{W^{2}}{w^{2}}+W^{2}\right)}} . \tag{4.89}
\end{equation*}
$$

The rest of (4.83) is the correlation term $\delta \rho(\varepsilon, \xi, \Delta)$ to this "naive" population which actually plays the dominant role in the long-time limit. As earlier, to calculate this term, we introduce the variables $\zeta=\xi-\varepsilon, \eta=(\varepsilon+\xi) / 2$ and write

$$
\begin{align*}
\delta \rho(\zeta, \eta, \Delta)= & \frac{1}{\left(\eta-\frac{\zeta}{2}-\Delta_{m}\right)\left(\eta-\frac{\zeta}{2}+i W\right)-\frac{W^{2}}{w^{2}}} \\
& \frac{1}{\left(\eta+\frac{\zeta}{2}-\Delta_{m}\right)\left(\eta+\frac{\zeta}{2}-i W\right)-\frac{W^{2}}{w^{2}}} \\
& \frac{-2 i W^{5} / N w^{4}}{\left[\left(\eta-\frac{\zeta}{2}+i W\right)\left(\eta+\frac{\zeta}{2}-i W\right)+\frac{W^{2}}{w^{2}}\right] \zeta}
\end{align*}
$$

The first and second factors have singularities at the points $\eta=\zeta / 2+$ $\varepsilon_{1,2}$ and $\eta=-\zeta / 2+\varepsilon_{1,2}^{*}$, whereas the third factor has poles at $\eta=$ $\pm \sqrt{(\zeta / 2-i W)^{2}-W^{2} / w^{2}}$. It is sufficient to consider only one of the first four singularities, say $\varepsilon_{1}$, since all the other poles of the first two parts result from $\varepsilon_{1}$ after complex conjugation associated with the replacement $\zeta \rightarrow-\zeta$, and after the replacement $\Delta \rightarrow-\Delta$. One has to take the real and even in $\Delta$ part of the final expression at the end of the calculations. In the third factor we retain only one pole which is located in the part of the complex plane which is opposite with respect to the real axis to the part where $\varepsilon_{1}$ is located, or else the integral over $d \eta$ vanishes. This results in

$$
\begin{align*}
& \delta \rho(\zeta, \eta, \Delta)=\frac{1}{\left(\eta-\frac{\zeta}{2}-\varepsilon_{1}\right)\left(\zeta+\varepsilon_{1}-\varepsilon_{1}^{*}\right)\left(\zeta+\varepsilon_{1}-\varepsilon_{2}^{*}\right)\left(\varepsilon_{1}-\varepsilon_{2}\right)} \\
& \frac{1}{\sqrt{\left(\frac{\zeta}{2}-i W\right)^{2}-\frac{W^{2}}{w^{2}}}} \frac{i W^{5} / N w^{4}}{\left[\left(\eta+\sqrt{\left(\frac{\zeta}{2}-i W\right)^{2}-\frac{W^{2}}{w^{2}}}\right)\right] \zeta} \tag{4.91}
\end{align*}
$$

and yields after integration over $d \eta$

$$
\begin{align*}
\delta \rho(\zeta, \Delta)= & \frac{-1}{\left(\sqrt{\left(\frac{\zeta}{2}-i W\right)^{2}-\frac{W^{2}}{w^{2}}}+\frac{\zeta}{2}+\varepsilon_{1}\right)\left(\varepsilon_{1}-\varepsilon_{2}\right) \zeta} \\
& \frac{2 \pi W^{5} / N w^{4}}{\left(\zeta+\varepsilon_{1}-\varepsilon_{1}^{*}\right)\left(\zeta+\varepsilon_{1}-\varepsilon_{2}^{*}\right) \sqrt{\left(\frac{\zeta}{2}-i W\right)^{2}-\frac{W^{2}}{w^{2}}}},
\end{align*}
$$

which is equivalent to an expression more convienient for further calculations

$$
\begin{align*}
\delta \rho(\zeta, \Delta)= & \frac{\left(\sqrt{(\zeta / 2-i W)^{2}-W^{2} / w^{2}}-\zeta / 2-\varepsilon_{1}\right)}{\left(-\left(\frac{\zeta}{2}-i W\right)^{2}+\frac{W^{2}}{w^{2}}+\left(\frac{\zeta}{2}+\varepsilon_{1}\right)^{2}\right)\left(\varepsilon_{1}-\varepsilon_{2}\right) \zeta} \\
& \frac{2 \pi W^{5} / N w^{4}}{\left(\zeta+\varepsilon_{1}-\varepsilon_{1}^{*}\right)\left(\zeta+\varepsilon_{1}-\varepsilon_{2}^{*}\right) \sqrt{\left(\frac{\zeta}{2}-i W\right)^{2}-\frac{W^{2}}{w^{2}}}} \tag{4.93}
\end{align*}
$$

One sees clearly that (4.92) has poles at the points $\zeta_{1,2,3}=0, \varepsilon_{1}^{*}-\varepsilon_{1}, \varepsilon_{2}^{*}-\varepsilon_{1}$, and at the point $\zeta_{4}=\left(W^{2}+W^{2} / w^{2}+\varepsilon_{1}^{2}\right) /\left(i W-\varepsilon_{1}\right)$ corresponding to the first factor in the denominator. Apart from these singularities, the expression has branching points at $\zeta_{5,6}=i W \pm W / w$. We note that the part

$$
\begin{equation*}
\delta \rho_{1}(\zeta, \eta, \Delta)=\frac{\left(i W-\varepsilon_{1}\right) 2 \pi W^{5} / N w^{4}\left(\varepsilon_{1}-\varepsilon_{2}\right)}{\zeta\left(\zeta-\zeta_{2}\right)\left(\zeta-\zeta_{3}\right)\left(\zeta-\zeta_{4}\right)} \tag{4.94}
\end{equation*}
$$

of the right-hand side of (4.93) does not contain the branching points, and after the inverse Fourier transform over $\zeta$ it gives a constant and three exponential time dependencies, whereas only the remainder

$$
\begin{equation*}
\delta \rho_{2}(\zeta, \Delta)=\frac{\left(-\zeta / 2-\varepsilon_{1}\right)\left(i W-\varepsilon_{1}\right) 4 \pi W^{5} / N w^{4}\left(\varepsilon_{1}-\varepsilon_{2}\right)}{\zeta\left(\zeta-\zeta_{2}\right)\left(\zeta-\zeta_{3}\right)\left(\zeta-\zeta_{4}\right) \sqrt{\left(\zeta-\zeta_{5}\right)\left(\zeta-\zeta_{6}\right)}} \tag{4.95}
\end{equation*}
$$

which does contain the branching points requires further consideration.
By casting the polynomial fraction factor in (4.95) as a sum of poles

$$
\begin{align*}
& \frac{\left(-\zeta / 2-\varepsilon_{1}\right)}{\zeta\left(\zeta-\zeta_{2}\right)\left(\zeta-\zeta_{3}\right)\left(\zeta-\zeta_{4}\right)}=\frac{\varepsilon_{1}}{\zeta_{2} \zeta_{3} \zeta_{4}} \frac{1}{\zeta}-\frac{\zeta_{2} / 2+\varepsilon_{1}}{\zeta_{2}\left(\zeta_{2}-\zeta_{3}\right)\left(\zeta_{2}-\zeta_{4}\right)} \frac{1}{\left(\zeta-\zeta_{2}\right)} \\
& -\frac{\zeta_{3} / 2+\varepsilon_{1}}{\zeta_{3}\left(\zeta_{3}-\zeta_{2}\right)\left(\zeta_{3}-\zeta_{4}\right)} \frac{1}{\left(\zeta-\zeta_{3}\right)}-\frac{\zeta_{4} / 2+\varepsilon_{1}}{\zeta_{4}\left(\zeta_{4}-\zeta_{2}\right)\left(\zeta_{4}-\zeta_{3}\right)} \frac{1}{\left(\zeta-\zeta_{4}\right)} \tag{4.96}
\end{align*}
$$

one reduces (4.95) to the form

$$
\begin{equation*}
\delta \rho_{2}(\zeta, \Delta)=\sum_{k} a_{k} \frac{\left(i W-\varepsilon_{1}\right) 4 \pi W^{5} / N w^{4}\left(\varepsilon_{1}-\varepsilon_{2}\right)}{\left(\zeta-\zeta_{k}\right) \sqrt{\left(\zeta-\zeta_{5}\right)\left(\zeta-\zeta_{6}\right)}} \tag{4.97}
\end{equation*}
$$

where the coefficients $a_{k}$ emerge from (4.96). The next step is the calculation of the inverse Fourier transform of the expression of type

$$
\begin{equation*}
T_{F}\left(t, \zeta_{k}, \zeta_{5}, \zeta_{6}\right)=\frac{1}{4 \pi^{2}} \int_{-\infty}^{\infty} \frac{\exp [i t \zeta]}{\left(\zeta-\zeta_{k}\right) \sqrt{\left(\zeta-\zeta_{5}\right)\left(\zeta-\zeta_{6}\right)}} d \zeta \tag{4.98}
\end{equation*}
$$

which is easier to perform by substituting explicitly $\zeta_{5,6}=i W \pm W / w$. After the replacement $\zeta \rightarrow x W / w+i W$, equation (4.98) takes the form

$$
\begin{equation*}
T_{F}\left(t, \zeta_{k}, \zeta_{5}, \zeta_{6}\right)=\frac{1}{4 \pi^{2}} \int_{-\infty}^{\infty} \frac{\exp [i t(x W / w+i W)]}{\left(x W / w+i W-\zeta_{k}\right) \sqrt{x^{2}-1}} d x \tag{4.99}
\end{equation*}
$$

We now multiply both parts of this equation by $\exp \left(-\zeta_{k} i t\right)$, take the time derivative, and arrive at

$$
\begin{align*}
\frac{d\left[e^{-i t \zeta_{k}} T_{F}\left(t, \zeta_{k}\right)\right]}{d t} & =\int_{-\infty}^{\infty} \frac{\exp \left[i t\left(x \frac{W}{w}+i W-\zeta_{k}\right)\right]}{\left(x \frac{W}{w}+i W-\zeta_{k}\right) \sqrt{x^{2}-1}} \frac{d x}{4 \pi^{2}}  \tag{4.100}\\
& =\frac{e^{-t\left(W+i \zeta_{k}\right)}}{4 \pi^{2}} \int_{-\infty}^{\infty} \frac{\exp \left[i t x \frac{W}{w}\right]}{\sqrt{x^{2}-1}} d x=e^{-t\left(W+i \zeta_{k}\right)} \frac{J_{0}(t W / w)}{2 \pi}
\end{align*}
$$

where we have made use of the same reasoning as for (4.78), (4.79). This implies that

$$
\begin{equation*}
T_{F}\left(t, \zeta_{k}\right)=e^{i t \zeta_{k}} \int_{0}^{t} e^{-\tau\left(W+i \zeta_{k}\right)} J_{0}\left(\tau \frac{W}{w}\right) \frac{d \tau}{2 \pi} \tag{4.101}
\end{equation*}
$$

where the lower limit is set to 0 , since at $t=0$ the integration contour (4.99) can be moved to infinity, where the integral vanishes.

Equations (4.94)-(4.98) with allowance for (4.89), (4.101) yield a clumsy albeit explicit expression for the distribution of the population over the band as a function of time. The result of such a calculation for $w=1 / \sqrt{3}$ is shown in Fig. 4.21(b). One sees an important feature of the distribution - it manifests oscillations in time and oscillations as a function of detuning until the asymptotic Lorentzian profile is attained. The larger the ratio $1 / w=$ $\sqrt{\left\langle V^{2}\right\rangle N} / \pi g\left\langle V^{2}\right\rangle$, the more pronounced the oscillations.

### 4.4.3 Role of the Interaction Rank

Some important aspects of the degenerate level-band problem, including the difference of the factor of 2 in the extreme $\sqrt{\left\langle V^{2}\right\rangle N} \gg \pi g\left\langle V^{2}\right\rangle$ between the exact asymptotic distribution of the populations and that suggested by the equal distribution principle, are revealed when we consider this system as an ensemble of two-level systems, by analogy with the consideration of two degenerate levels presented in Sect.4.3.2. As earlier, in (4.62) we concentrate here on the matrix structure of the Hamiltonian

$$
\hat{H}=\left(\begin{array}{cc}
\hat{H}_{2} & \hat{V}  \tag{4.102}\\
\hat{V}^{+} & \hat{H}_{1}
\end{array}\right)=\left(\begin{array}{cc}
\Delta_{m} \delta_{m}^{m^{\prime}} & V_{n}^{m} \\
V_{m}^{n} & 0
\end{array}\right)
$$

and truncate the band such that only the levels that were populated in the process of population transfer are taken into account. The Hamiltonian now also consists of four blocks, a small $M \times M$ block $\hat{H}_{2}$ of order $M \sim$ $\pi \sqrt{\left\langle V^{2}\right\rangle N g^{2}}$, a large $N \times N$ block $\hat{H}_{1}$ of order $N$, and two $M \times N$ and $N \times M$ blocks $V_{n}^{m}$ and $V_{m}^{n}$ of the interaction matrices, that apparently have a rank not larger than $M$. Since $\hat{H}_{1}=0$, one can perform a linear unitary transformation such that the rectangular conjugated matrices $\hat{V}$ and $\hat{V}^{+}$have only non-zero upper and right $M \times M$ blocks, while all of the other matrix elements vanish. The Hamiltonian (4.102) then takes the form

$$
\hat{H}=\left(\begin{array}{ccc}
\hat{H}_{2} & \hat{V}^{\prime} & 0  \tag{4.103}\\
\hat{V}^{\prime}+ & 0 & 0 \\
0 & 0 & 0
\end{array}\right)=\left(\begin{array}{cc}
\Delta_{m} \delta_{m}^{m^{\prime}} & V_{n}^{m} \\
V_{m}^{n} & 0
\end{array}\right)
$$

where $\hat{V}^{\prime}$ and $\hat{V}^{\prime}+$ are hermitian conjugate $M \times M$ matrices of rank $M^{\prime}$ not higher than $M$. In other words, via the linear unitary transformation applied
to the states of the degenerate level we arrive at a representation where only the first $M^{\prime}$ states of the basis are coupled to the upper band. These states are composed of the properly normalized and orthogonalized linear combinations of $M^{\prime}$ different states

$$
\begin{equation*}
\left|c_{m}\right\rangle=\sum_{n=0}^{N} V_{m}^{n}|n\rangle \tag{4.104}
\end{equation*}
$$

interacting with the band, whereas the remaining $N-M^{\prime}$ states are orthogonal to these $M^{\prime}$ states. One calls the latter "dark states", as opposed to "bright" states (4.104), since they do not interact with the band. The completely decoupled dark states retain their initial population, and hence the excitation process leaves the $\left(N-M^{\prime}\right) / N$ fraction of the initial population of the degenerate level intact. The rest of the population $M^{\prime} / N$ is equally distributed among $M$ states of band 2 and $M^{\prime}$ bright states of the degenerate level, provided the matrix $\hat{V}^{\prime}$ commutes with its Hermitian conjugate $\hat{V}^{\prime}+$ and can be diagonalized by a unitary transformation, which we assume to be the case. For $M^{\prime}=M$ this yields $\rho_{2}=M / 2 N, \rho_{1}=(N-M) / N+M / 2 N$, and we arrive at

$$
\begin{equation*}
\rho_{2} / \rho_{1} \simeq M / 2 N=\pi\left\langle V^{2}\right\rangle g / \sqrt{\left\langle V^{2}\right\rangle N} \tag{4.105}
\end{equation*}
$$

which coincides with the extreme following from (4.82) and differs by the factor of 2 from the estimate on p. 156

In the example just considered we have seen that the rank of the interaction matrix, that is the number of its non-zero eigenvalues, may play an important role in the population dynamics of a complex multilevel quantum system. It effects the number of states that actually participate in the process. The rank of the interaction matrix is closely related to the question of the applicability of the random matrix model to the description of quantum dynamics in a complex system. Indeed, a square $N \times N$ matrix which has rank $R \ll N$ can be cast with the help of the Schmidt expansion into the sum

$$
\begin{equation*}
V_{j}^{k}=\sum_{s=1}^{R}\left\langle j \mid c_{s}\right\rangle V_{s}\left\langle c_{s} \mid k\right\rangle \tag{4.106}
\end{equation*}
$$

over its eigenvalues $V_{s}$ and eigenvectors $\left|c_{s}\right\rangle$. In the extreme case $R=1$ this yields

$$
\begin{equation*}
V_{j}^{k}=\left\langle j \mid c_{1}\right\rangle V_{1}\left\langle c_{1} \mid k\right\rangle=C_{j} V_{1} C_{k}^{*} \tag{4.107}
\end{equation*}
$$

with $C_{j}=\left\langle j \mid c_{1}\right\rangle$, which means that the interaction matrix assumes a factorized form. Note that the matrix elements of $\hat{V}$ can still be random, if the expansion coefficients $\left\langle c_{1} \mid k\right\rangle$ of its eigenvector $\left|c_{1}\right\rangle$ over the basis set of the energy eigenstates $|k\rangle$ of the unperturbed Hamiltonian are random. In this case, however, as well as in the more general case $R \ll N$, the statistical properties of the products of the matrix elements differ drastically from that for completely random matrices, and in particular this circumstance does not
allow one to select from the diagram series of Fig. 4.5 the "tree"-like contributions of Fig. 4.6. This implies a strong statistical correlation among the matrix elements of $\hat{V}$, which one can see explicitly for (4.107) for which we have

$$
\begin{equation*}
\left.\left.\langle | V_{j}^{k} V_{k}^{r}\right|^{2}\right\rangle=\left(V_{1}\right)^{4} / N^{2} \tag{4.108}
\end{equation*}
$$

and not

$$
\begin{equation*}
\left\langle V_{j}^{k} V_{k}^{r}\right\rangle=\left\langle V^{2}\right\rangle \delta_{j}^{r} \tag{4.109}
\end{equation*}
$$

as would be the case for completely random matrices.

### 4.5 The Role of Correlations

The correlations among the interaction matrix elements allow for the interference of different excitation trajectories which not only do not disappear after the ensemble average, but on the contrary strongly effect the population dynamics in the quantum system in general, and the population distribution among the states of different energies in particular. We illustrate this with two examples, that of two levels interacting via a band, and a two-band system with an interaction matrix of rank 1 given in (4.107).

### 4.5.1 Two levels and a band

Let us consider a quantum system of two levels, each of which interacts with a band of levels as shown in Fig. 4.22. This system resembles the Fano system


Fig. 4.22. Two levels coupled to a band by different interaction matrix elements $V_{0}^{n}$, and $V_{1}^{n}$.
depicted in Fig. 3.7, that was previously considered in Sect.3.2.2, although here we will not assume equal couplings $V$ and $V^{\prime}$ of each of the levels to the states of the band, but on the contrary we concentrate on the statistical properties of these matrix elements.

## The General Solution

By analogy to (3.49) we can write

$$
\begin{align*}
& \varepsilon \psi_{0}=\sum_{n=1}^{N} V_{0}^{n} \psi_{n}+i \\
& \varepsilon \psi_{1}=\sum_{n=1}^{N} V_{1}^{n} \psi_{n}+\Delta_{1} \psi_{1}+V^{\prime \prime} \psi_{0} \\
& \varepsilon \psi_{n}=\Delta_{n} \psi_{n}+V_{n}^{0} \psi_{0}+V_{n}^{1} \psi_{1} . \tag{4.110}
\end{align*}
$$

where $\psi_{0}(t=0)=1 ; \psi_{1}(t=0)=0 ; \psi_{n}(t=0)=0$ are taken as the initial conditions. Substitution of $\psi_{n}$ from the last equation into the first two equations yields the system of equations

$$
\begin{align*}
\left(\varepsilon-\sum_{n=1}^{N} \frac{V_{0}^{n} V_{n}^{0}}{\left(\varepsilon-\Delta_{n}\right)}\right) \psi_{0}-\sum_{n=1}^{N} \frac{V_{0}^{n} V_{n}^{1}}{\left(\varepsilon-\Delta_{n}\right)} \psi_{1} & =i \\
-\sum_{n=1}^{N} \frac{V_{1}^{n} V_{n}^{0}}{\left(\varepsilon-\Delta_{n}\right)} \psi_{0}+\left(\varepsilon-\sum_{n=1}^{N} \frac{V_{1}^{n} V_{n}^{1}}{\left(\varepsilon-\Delta_{n}\right)}-\Delta_{1}\right) \psi_{1} & =0 . \tag{4.111}
\end{align*}
$$

which have the solutions

$$
\begin{align*}
& \psi_{0}(\varepsilon)=\frac{i\left(\varepsilon-W_{1}^{1}-\Delta_{1}\right)}{\left(\varepsilon-W_{0}^{0}\right)\left(\varepsilon-W_{1}^{1}-\Delta_{1}\right)-W_{0}^{1} W_{1}^{0}} \\
& \psi_{1}(\varepsilon)=\frac{i W_{1}^{0}}{\left(\varepsilon-W_{0}^{0}\right)\left(\varepsilon-W_{1}^{1}-\Delta_{1}\right)-W_{0}^{1} W_{1}^{0}} \tag{4.112}
\end{align*}
$$

where we denote $W_{p}^{q}=\sum_{n=1}^{N} V_{p}^{n} V_{n}^{q} /\left(\varepsilon-\Delta_{n}\right)$ with $p, q=0,1$.
Let us now choose the interactions in the form

$$
\begin{align*}
V_{n}^{0} & =V C_{n} \\
V_{n}^{1} & =V S_{n} . \tag{4.113}
\end{align*}
$$

similar to (4.107), where $V$ is a coupling constant of energy dimensionality whereas the coefficients $C_{n}=\left|C_{n}\right| e^{i \theta_{n}}$ and $S_{n}=\left|S_{n}\right| e^{i \chi_{n}}$ are complex numbers normalized by the conditions $\sum_{n}\left|C_{n}\right|^{2}=\sum_{n}\left|S_{n}\right|^{2}=1$. Equation (4.113) implies that the state $|0\rangle$ interacts with only one linear combination $|c\rangle=\sum_{n} C_{n}|n\rangle$ of the states of the band and that for the state $|1\rangle$ the similar combination $|s\rangle=\sum_{n} S_{n}|n\rangle$ of the band levels might be different. Here we assume equal interaction matrix elements, that is $\sum_{n}\left|V_{0}^{n}\right|^{2}=\sum_{n}\left|V_{1}^{n}\right|^{2}$, which is actually not a restricting assumption, since the strip of interacting levels can have an arbitrary width.

Substitution of (4.113) allows one to express the coefficients $W_{p}^{q}$ in (4.112) in the form

$$
\begin{align*}
& W_{0}^{0}=V^{2} \mathcal{S}_{1}=V^{2} \sum_{n=1}^{N}\left|C_{n}\right|^{2} /\left(\varepsilon-\Delta_{n}\right), \\
& W_{1}^{1}=V^{2} \mathcal{S}_{2}=V^{2} \sum_{n=1}^{N}\left|S_{n}\right|^{2} /\left(\varepsilon-\Delta_{n}\right), \\
& W_{1}^{0}=\left(W_{0}^{1}\right)^{*}=V^{2} \mathcal{S}_{3}=V^{2} \sum_{n=1}^{N} C_{n}^{*} S_{n} /\left(\varepsilon-\Delta_{n}\right) . \tag{4.114}
\end{align*}
$$

We now assume that both $\left|C_{n}\right|^{2}$ and $\left|S_{n}\right|^{2}$ are smooth functions of the index $n$, and in the vicinity of $\Delta_{n}=0$ they amount to $|C|^{2}$ and $|S|^{2}$ respectively. For the regime $t<g$ we can therefore factor out these quantities from the sums (4.114) and replace the first two sums by integrations, which yields

$$
\begin{align*}
& \mathcal{S}_{1}=\sum_{n=1}^{N} \frac{\left|C_{n}\right|^{2}}{\varepsilon-\Delta_{n}}=i \pi g|C|^{2}, \\
& \mathcal{S}_{2}=\sum_{n=1}^{N} \frac{\left|S_{n}\right|^{2}}{\varepsilon-\Delta_{n}}=i \pi g|S|^{2} . \tag{4.115}
\end{align*}
$$

The sum $\mathcal{S}_{3}=V^{2} \sum_{n=1}^{N} C_{n}^{*} S_{n} /\left(\varepsilon-\Delta_{n}\right)$ requires a more detailed analysis. Indeed, apart from the moduli $\left|C_{n}\right|$ and $\left|S_{n}\right|$ it contains the phase factors $e^{i \theta_{n}-i \chi_{n}}$, that may differ considerably from unity, thus adding another dependence of the summand on the summation index $n$. However, for $t \ll g$ this problem can be considerably simplified when we make use of the uncertainty principle. The simplest way to employ this principle it is to say that at a time $t$, the typical frequency $\varepsilon$ bringing the main contribution to the integral of the inverse Fourier transformation is of the order of $1 / t$, and that generally speaking, this number has a complex value. Denominators of the sums (4.114) do not therefore differ much for two sequential terms, since the difference $\Delta_{n+1}-\Delta_{n} \sim 1 / g$ remains less than the imaginary part of $\varepsilon \sim 1 / t$ in the denominator governing the resonance width by the condition $t \ll g$. We can therefore perform the summation for $\mathcal{S}_{3}$ in two steps: by taking first the average of the phase factors over a number $2 K \sim g / t$ of neighboring terms for which the denominators $\left(\varepsilon-\Delta_{n}\right)$ are almost identical, and then by replacing these factors by their average $\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle=\sum_{m=n-K}^{n+K} e^{i \theta_{m}-i \chi_{m}} / 2 K$, which is likely a smooth function of the index $n$. Therefore replacing summation by integration one arrives at

$$
\begin{equation*}
\mathcal{S}_{3}=V^{2} \sum_{n=1}^{N} \frac{C_{n}^{*} S_{n}}{\varepsilon-\Delta_{n}}=i \pi g|C S|\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle \tag{4.116}
\end{equation*}
$$

where, in the last part, the index $n$ corresponds to the states for which $\Delta_{n} \sim$ 0 . By taking such an average we assume implicitly that the absolute values of $C_{n}$ and $S_{n}$ and their phases are statistically independent.

Substitution of (4.114)-(4.116) into (4.112) yields

$$
\begin{align*}
& \psi_{0}(\varepsilon)=i \frac{\varepsilon-\Delta_{1}+i \pi g V^{2}|C|^{2}}{\left(\varepsilon-\varepsilon_{1}\right)\left(\varepsilon-\varepsilon_{2}\right)} \\
& \psi_{1}(\varepsilon)=-\frac{\pi g V^{2}|C S|\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle}{\left(\varepsilon-\varepsilon_{1}\right)\left(\varepsilon-\varepsilon_{2}\right)} \tag{4.117}
\end{align*}
$$

where

$$
\begin{align*}
\varepsilon_{1,2}= & \frac{\Delta_{1}+g i \pi V^{2}\left(|C|^{2}+|S|^{2}\right)}{2} \pm\left[\frac{\left(\Delta_{1}+g i \pi V^{2}\left(|C|^{2}+|S|^{2}\right)\right)^{2}}{4}\right. \\
& \left.-g \pi V^{2}|C|^{2}\left(i \Delta_{1}+\left(\left|\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle\right|^{2}-1\right) g \pi V^{2}|S|^{2}\right)\right]^{1 / 2} \tag{4.118}
\end{align*}
$$

are the roots of the denominators responsible for the time evolution of the system at $t<g$. One sees that the phase correlation of the matrix elements $\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle$ is an important ingredient in these expressions.

## Different Particular Cases

In the particular case of totally random phases, $\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle=0$, and hence, according to (4.118) the probability amplitude of being in state $|1\rangle$ vanishes completely. For this case $\varepsilon_{1}=\Delta_{1}-i \pi g V^{2}|C|^{2}$, and $\varepsilon_{2}=-i \pi g V^{2}|S|^{2}$, which yields the result

$$
\begin{equation*}
\psi_{0}(t)=\frac{1}{2 \pi} \int_{C} \frac{i e^{-i \varepsilon t} d \varepsilon}{\left(\varepsilon+i \pi g V^{2}|S|^{2}\right)}=e^{-W_{0}^{0} t} \tag{4.119}
\end{equation*}
$$

coinciding with the corresponding expressions of (3.32)-(3.33). This means that in the absence of correlation among the matrix elements that couple the states $|0\rangle$ and $|1\rangle$ to the band, each of these states interacts with the states of the band independently, thus conforming to the level-band model and yielding all the consequences of this model including (3.35) for the distribution of the population over the band levels. Therefore the state $|1\rangle$, which was not populated initially, remains unpopulated as long as $t<g$, and the continuous-band model remains valid, whereas the population of state $|0\rangle$ prefers to be distributed amongst the more statistically abundant states of the band.

On the other hand, in the case of strongest correlation $\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle=1$ the most pronounced effect takes place for the case $\Delta_{1}=0$, when

$$
\varepsilon_{1,2}=\left\{\begin{array}{c}
0  \tag{4.120}\\
-i \pi g V^{2}\left(|C|^{2}+|S|^{2}\right)
\end{array}\right.
$$

and therefore (4.117) take the form

$$
\begin{align*}
& \psi_{0}(\varepsilon)=i \frac{\varepsilon+i \pi g V^{2}|C|^{2}}{\varepsilon\left(\varepsilon+i \pi g V^{2}\left(|C|^{2}+|S|^{2}\right)\right)} \\
& \psi_{1}(\varepsilon)=i \frac{i \pi g V^{2}|C S|}{\varepsilon\left(\varepsilon+i \pi g V^{2}\left(|C|^{2}+|S|^{2}\right)\right)} \tag{4.121}
\end{align*}
$$

and yield

$$
\begin{align*}
& \psi_{0}(t)=\frac{|C|^{2}}{|C|^{2}+|S|^{2}}+\frac{|S|^{2}}{|C|^{2}+|S|^{2}} e^{-t \pi g V^{2}\left(|C|^{2}+|S|^{2}\right)} \\
& \psi_{1}(t)=\frac{|C S|}{|C|^{2}+|S|^{2}}\left(1-e^{-t \pi g V^{2}\left(|C|^{2}+|S|^{2}\right)}\right) \tag{4.122}
\end{align*}
$$

The most important consequence of the complete correlation of the matrix elements is that $\varepsilon_{1}=0$, and hence the probability amplitudes of (4.122) do not vanish as $t \rightarrow \infty$. The levels retain a significant part of the total population $\rho_{0}(t=\infty)=|C|^{4} /\left(|C|^{2}+|S|^{2}\right)^{2}, \rho_{1}(t=\infty)=|C|^{2}|S|^{2} /\left(|C|^{2}+\right.$ $\left.|S|^{2}\right)^{2}$, whereas the total population of the band levels amounts to $\rho_{b}(t=$ $\infty)=\left(|S|^{4}+|C|^{2}|S|^{2}\right) /\left(|C|^{2}+|S|^{2}\right)^{2}$. In the particular case $|C|^{2}=|S|^{2}=1 / 2$ both of the levels have the asymptotic population $\rho_{0}(t=\infty)=\rho_{1}(t=\infty)=$ $1 / 4$, while the band takes the other half of the population.

The last case of $|C|^{2}=|S|^{2}=1 / 2,\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle=1$ and $\Delta_{1}=0$ can be readily understood when we consider the population dynamics of the system taking for the basis set of the problem the combinations $|+\rangle=(|0\rangle+|1\rangle) / \sqrt{2}$, and $|-\rangle=(|0\rangle-|1\rangle) / \sqrt{2}$, that is the symmetric and antisymmetric combinations of the levels. The state $|+\rangle$ interacts with the band states via a constant matrix element $V$ whereas the state $|-\rangle$ turns out to be completely decoupled. Therefore the state $|-\rangle$ retains the initial population $|\langle 1 \mid-\rangle|^{2}=1 / 2$ and the other half of the population, which was initially in the state $|+\rangle$, moves completely to the levels of the band. We note that in the case where $\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle=a<1$, the combination $|-\rangle$ starts to decay as well, since

$$
\begin{equation*}
\varepsilon_{1,2}=\frac{1}{2} g i \pi V^{2} \pm \frac{1}{2} g i \pi V^{2}|a| \tag{4.123}
\end{equation*}
$$

which yields the amplitudes

$$
\begin{align*}
& \psi_{+}(t)=\frac{1}{\sqrt{2}} e^{-t \pi g V^{2}(1+|a|) / 2} \\
& \psi_{-}(t)=\frac{1}{\sqrt{2}} e^{-t \pi g V^{2}(1-|a|) / 2} \tag{4.124}
\end{align*}
$$

and the populations

$$
\begin{align*}
& \rho_{0}(t)=\frac{1}{4}\left(e^{-t \pi g V^{2}(1+|a|) / 2}+e^{-t \pi g V^{2}(1-|a|) / 2}\right)^{2} \\
& \rho_{1}(t)=\frac{1}{4}\left(e^{-t \pi g V^{2}(1+|a|) / 2}-e^{-t \pi g V^{2}(1-|a|) / 2}\right)^{2} \tag{4.125}
\end{align*}
$$

The last terms manifest a tri-exponential decay with rates $\pi g V^{2}(1+|a|)$, $\pi g V^{2}$, and $\pi g V^{2}(1-|a|)$.

One more simple limiting case with $|C|^{2}=|S|^{2}$ (we take both equal to $1 / 2),\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle=1$ and $\Delta_{1} \neq 0$ corresponds to two resonant levels interacting via a continuum. Equation (4.118) yields for this case

$$
\begin{equation*}
\varepsilon_{1,2}=\frac{1}{2}\left[\Delta_{1}+g i \pi V^{2}\right] \pm \frac{1}{2}\left[\Delta_{1}^{2}-\left(g \pi V^{2}\right)^{2}\right]^{1 / 2} \tag{4.126}
\end{equation*}
$$

which after substitution into (4.117) and the inverse Fourier transformation results in

$$
\begin{align*}
\psi_{0}(t)= & \frac{-i \Delta_{1} e^{i t\left[\Delta_{1}+g i \pi V^{2}\right] / 2}}{\sqrt{\Delta_{1}^{2}-\left(g \pi V^{2}\right)^{2}}} \sin \left(\frac{t}{2} \sqrt{\Delta_{1}^{2}-\left(g \pi V^{2}\right)^{2}}\right) \\
& +e^{i t\left[\Delta_{1}+g i \pi V^{2}\right] / 2} \cos \left(\frac{t}{2} \sqrt{\Delta_{1}^{2}-\left(g \pi V^{2}\right)^{2}}\right) \\
\psi_{1}(\varepsilon)= & \frac{g \pi V^{2} e^{i t\left[\Delta_{1}+g i \pi V^{2}\right] / 2}}{\sqrt{\Delta_{1}^{2}-\left(g \pi V^{2}\right)^{2}}} \sin \left(\frac{t}{2} \sqrt{\Delta_{1}^{2}-\left(g \pi V^{2}\right)^{2}}\right) . \tag{4.127}
\end{align*}
$$

From this equation one sees that the population

$$
\begin{equation*}
\rho_{1}(t)=\frac{\left(g \pi V^{2}\right)^{2}}{\Delta_{1}^{2}-\left(g \pi V^{2}\right)^{2}} e^{-t g \pi V^{2}} \sin ^{2}\left(\frac{t}{2} \sqrt{\Delta_{1}^{2}-\left(g \pi V^{2}\right)^{2}}\right) \tag{4.128}
\end{equation*}
$$

of state $|1\rangle$ manifests two different dynamical regimes of dying oscillations for $\left|\Delta_{1}\right|>\left|g \pi V^{2}\right|$, and a bi-exponential decay in the opposite case. For a fixed $\Delta_{1}$ the decay rate depends on $V$. For large $V$, when $\sqrt{\left(g \pi V^{2}\right)^{2}-\Delta_{1}^{2}} \simeq$ $g \pi V^{2}-\Delta_{1}^{2} / 2 g \pi V^{2}$ and the slowest exponent gives the main contribution to the population, (4.128) reads

$$
\begin{equation*}
\rho_{1}(t)=\frac{1}{4} e^{-t \Delta_{1}^{2} / 2 g \pi V^{2}} . \tag{4.129}
\end{equation*}
$$

One sees that the slowest decay rate $\Delta_{1}^{2} / 2 g \pi V^{2}$ decreases with the increase of the interaction $V$. Insight into this paradoxical result is gained by noting that the states $|0\rangle$ and $|1\rangle$ are coupled via the band, and that the stronger the interaction is, the larger the coupling. With an extremely strong coupling, the eigenstates of a two-level system (3.21) are very close to the combinations
 the case $\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle=1$, as shown above.

## Distribution of the Population over the Band

From the analysis of the two-level-band system we have learned that the correlations of the interaction matrix elements may have an important influence
on the population dynamics. Interference phenomena may result in a considerable increase in the decay time. When earlier, in Sect.3.2.2, considering the Fano problem, we saw that the interference phenomenon in a two-level band system can manifest titself also in the population distribution over the band. Now we will consider this phenomenon in more detail for the system of two levels and a band. We assume that $\varepsilon_{1}$ of (4.118) differs from zero and has therefore a non-zero imaginary part, since otherwise the system is equivalent to a level-band system that has already been considered in Chap. 3.

From the third equation (4.110) with the allowance of (4.113) one obtains

$$
\begin{equation*}
\psi_{n}=\frac{V C_{n}}{\left(\varepsilon-\Delta_{n}\right)} \psi_{0}+\frac{V S_{n}}{\left(\varepsilon-\Delta_{n}\right)} \psi_{1} \tag{4.130}
\end{equation*}
$$

and substitution of $\psi_{0}$ and $\psi_{1}$ from (4.117) yields

$$
\begin{align*}
\psi_{n}= & i \frac{V C_{n}}{\left(\varepsilon-\Delta_{n}\right)} \frac{\varepsilon-\Delta_{1}+i \pi g V^{2}|C|^{2}}{\left(\varepsilon-\varepsilon_{1}\right)\left(\varepsilon-\varepsilon_{2}\right)} \\
& -\frac{V S_{n}}{\left(\varepsilon-\Delta_{n}\right)} \frac{\pi g V^{2}|C S|\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle}{\left(\varepsilon-\varepsilon_{1}\right)\left(\varepsilon-\varepsilon_{2}\right)} . \tag{4.131}
\end{align*}
$$

In the regime $1 / W_{m} \ll t \ll g$, the main contribution to the inverse Fourier transform comes from the pole $\varepsilon=\Delta_{n}$, since two other poles $\varepsilon=\varepsilon_{1,2}$ give exponentially small contributions that can be majorized by the exponent $e^{-t W_{m}}$, where $W_{m}$ is the minimum of two imaginary parts of $\varepsilon_{1,2}$. This yields

$$
\begin{align*}
\rho(\Delta)= & \frac{W^{\prime}}{\pi g}\left|\frac{1}{\left(\Delta-\varepsilon_{1}\right)\left(\Delta-\varepsilon_{2}\right)}\right|^{2} \\
& \left|C_{n}\left(\Delta-\Delta_{1}+i W^{\prime}|C|^{2}\right)+i S_{n} W^{\prime}\right| C S\left|\left\langle e^{i \theta_{n}-i \chi_{n}}\right\rangle\right|^{2} \tag{4.132}
\end{align*}
$$

for the population $\rho\left(\Delta_{n}\right)=\rho_{n}=\left|\psi_{n}\right|^{2}$ of a state $|n\rangle$, where by $W^{\prime}$ we denote the combination $\pi g V^{2}$.

By taking the average of the last factor in (4.132) over a small energy interval we obtain the mean population spectral density $g \rho(\Delta)$ in the form

$$
\begin{equation*}
g \rho(\Delta)=\frac{W u}{\pi} \frac{\left(\Delta-\Delta_{1}\right)^{2}+W^{2}\left(u^{2}+2 \alpha(1-u)\right)}{\left|\left(\Delta-\varepsilon_{1}\right)\left(\Delta-\varepsilon_{2}\right)\right|^{2}} \tag{4.133}
\end{equation*}
$$

where by $\alpha$ we denote $|a|^{2}=\left|\left\langle e^{i \chi_{n}-i \theta_{n}}\right\rangle\right|^{2}$, take $u=|C|^{2} /\left(|C|^{2}+|S|^{2}\right)$, $W=W^{\prime}\left(|C|^{2}+|S|^{2}\right)$, and write the roots

$$
\begin{equation*}
\frac{\varepsilon_{1,2}}{W}=\frac{\Delta_{1}}{2 W}+\frac{i}{2} \pm\left[\left(\frac{\Delta_{1}}{2 W}+\frac{i}{2}\right)^{2}-u\left(\frac{i \Delta_{1}}{W}+(\alpha-1)(1-u)\right)\right]^{1 / 2} \tag{4.134}
\end{equation*}
$$

in the same notation. The asymptotic population profile of (4.133) is a sum of two Lorentzian profiles

$$
\begin{equation*}
g \rho(\Delta)=\frac{1}{\pi} \frac{f_{1}}{\left|\Delta-\varepsilon_{1}\right|^{2}}+\frac{1}{\pi} \frac{f_{2}}{\left|\Delta-\varepsilon_{2}\right|^{2}}+\frac{f_{0}}{\pi}\left(\frac{\Delta}{\left|\Delta-\varepsilon_{1}\right|^{2}}-\frac{\Delta}{\left|\Delta-\varepsilon_{2}\right|^{2}}\right) \tag{4.135}
\end{equation*}
$$

that are placed at the positions $\operatorname{Re} \varepsilon_{1,2}$ and have the widths $\operatorname{Im} \varepsilon_{1,2}$ respectively and an asymmetric correction term. The corresponding strengths $f_{0,1,2}$ are algebraic functions of $\Delta_{1}, \alpha$, and $u$, that can be found by equating the coefficients of quadratic, linear and constant terms in the numerators of the right-hand sides of (4.133) and (4.135):

$$
\begin{align*}
& f_{1}+f_{2}-f_{0}\left(2 R_{2}-2 R_{1}\right)=W u \\
& R_{2} f_{1}+f_{2} R_{1}-f_{0} \frac{\left|\varepsilon_{2}\right|^{2}-\left|\varepsilon_{1}\right|^{2}}{2}=\Delta_{1} W u \\
& f_{1}\left|\varepsilon_{2}\right|^{2}+f_{2}\left|\varepsilon_{1}\right|^{2}=W u\left(\Delta_{1}^{2}+W^{2}\left(u^{2}+2 \alpha(1-u)\right)\right) \tag{4.136}
\end{align*}
$$

where $R_{1}=\operatorname{Re} \varepsilon_{1}, R_{2}=\operatorname{Re} \varepsilon_{2}$, which yields

$$
\begin{align*}
f_{0}= & \frac{2 u W\left(R_{2}-R_{1}\right)\left(\left(2 \alpha(1-u)+u^{2}\right) W^{2}+\Delta_{1}^{2}\right)}{4 R_{2}^{2}\left|\varepsilon_{1}^{2}\right|+4 R_{1}^{2}\left|\varepsilon_{2}^{2}\right|+\left(\left|\varepsilon_{1}^{2}\right|-\left|\varepsilon_{2}^{2}\right|\right)^{2}-4 R_{1} R_{2}\left(\left|\varepsilon_{1}^{2}\right|+\left|\varepsilon_{2}^{2}\right|\right)} \\
& +\frac{2 u W\left(R_{1}\left|\varepsilon_{2}^{2}\right|-R_{2}\left|\varepsilon_{1}^{2}\right|-\Delta_{1}\left(\left|\varepsilon_{2}^{2}\right|-\left|\varepsilon_{1}^{2}\right|\right)\right)}{4 R_{2}^{2}\left|\varepsilon_{1}^{2}\right|+4 R_{1}^{2}\left|\varepsilon_{2}^{2}\right|+\left(\left|\varepsilon_{1}^{2}\right|-\left|\varepsilon_{2}^{2}\right|\right)^{2}-4 R_{1} R_{2}\left(\left|\varepsilon_{1}^{2}\right|+\left|\varepsilon_{2}^{2}\right|\right)}, \\
f_{1}= & u W \frac{\left(\left(2 \alpha(1-u)+u^{2}\right) W^{2}+\Delta_{1}^{2}\right)\left(4 R_{1}\left(R_{1}-R_{2}\right)-\left|\varepsilon_{1}^{2}\right|+\left|\varepsilon_{2}^{2}\right|\right)}{4 R_{2}^{2}\left|\varepsilon_{1}^{2}\right|+4 R_{1}^{2}\left|\varepsilon_{2}^{2}\right|+\left(\left|\varepsilon_{1}^{2}\right|-\left|\varepsilon_{2}^{2}\right|\right)^{2}-4 R_{1} R_{2}\left(\left|\varepsilon_{1}^{2}\right|+\left|\varepsilon_{2}^{2}\right|\right)} \\
& +\frac{u W\left|\varepsilon_{1}^{2}\right|\left(4\left(R_{2}-R_{1}\right) \Delta_{1}+\left|\varepsilon_{1}^{2}\right|-\left|\varepsilon_{2}^{2}\right|\right)}{4 R_{2}^{2}\left|\varepsilon_{1}^{2}\right|+4 R_{1}^{2}\left|\varepsilon_{2}^{2}\right|+\left(\left|\varepsilon_{1}^{2}\right|-\left|\varepsilon_{2}^{2}\right|\right)^{2}-4 R_{1} R_{2}\left(\left|\varepsilon_{1}^{2}\right|+\left|\varepsilon_{2}^{2}\right|\right)}, \\
f_{2}= & u W \frac{\left(\left(2 \alpha(1-u)+u^{2}\right) W^{2}+\Delta_{1}^{2}\right)\left(4 R_{2}\left(R_{2}-R_{1}\right)+\left|\varepsilon_{1}^{2}\right|-\left|\varepsilon_{2}^{2}\right|\right)}{4 R_{2}^{2}\left|\varepsilon_{1}^{2}\right|+4 R_{1}^{2}\left|\varepsilon_{2}^{2}\right|+\left(\left|\varepsilon_{1}^{2}\right|-\left|\varepsilon_{2}^{2}\right|\right)^{2}-4 R_{1} R_{2}\left(\left|\varepsilon_{1}^{2}\right|+\left|\varepsilon_{2}^{2}\right|\right)} \\
& +\frac{u W\left|\varepsilon_{2}^{2}\right|\left(4\left(R_{1}-R_{2}\right) \Delta_{1}-\left|\varepsilon_{1}^{2}\right|+\left|\varepsilon_{2}^{2}\right|\right)}{4 R_{2}^{2}\left|\varepsilon_{1}^{2}\right|+4 R_{1}^{2}\left|\varepsilon_{2}^{2}\right|+\left(\left|\varepsilon_{1}^{2}\right|-\left|\varepsilon_{2}^{2}\right|\right)^{2}-4 R_{1} R_{2}\left(\left|\varepsilon_{1}^{2}\right|+\left|\varepsilon_{2}^{2}\right|\right)} . \tag{4.137}
\end{align*}
$$

One sees that the asymmetric correction term of the amplitude $f_{0}$ vanishes for identical positions of the poles $\varepsilon_{1}=\varepsilon_{2}$, but that in the general case it is of the order of two other terms that give the Lorentzian profiles. However, in the absence of direct interaction between the $|0\rangle$ and $|1\rangle$ levels, this correction term never compensates completely for the Lorentzian profiles, so that the population of the band never becomes zero at a band level with a finite detuning $\Delta_{m}$, which was the case for the Fano problem where such a direct coupling is present and renders the dependence of the coupling averaged over a small energy interval essentially non-uniform. One can see the resulting profiles in Fig. 4.23, where we show the dependence of the distribution $g \rho(\Delta)$ on the parameters $\Delta_{1}, \alpha$, and $u$. In principle, with the help of (4.136) one can find the correlation factor $\alpha$ of the band spectrum by analyzing the asymptotic population distribution $g \rho(\Delta)$ over the band states.


Fig. 4.23. The population distribution over a band as a function of the scaled detuning $\Delta / W$ for two levels coupled to a band by different interaction matrix elements $V_{0}^{n}=V C_{n}$ and $V_{1}^{n}=V S_{n}$. Parameters governing the distribution are (a) the detuning between the levels $\Delta_{1}$, divided by $\left.W=i \pi V^{2}\left(\left.\langle | C_{n}\right|^{2}\right\rangle+\left.\langle | S_{n}\right|^{2}\right\rangle$, (b) the relative coupling $\left.\left.u=\left.\langle | C_{n}\right|^{2}\right\rangle /\left(\left.\langle | C_{n}\right|^{2}\right\rangle+\left.\langle | S_{n}\right|^{2}\right\rangle$ of the level 0 , and (c) the correlation $\left.a=\left|\left\langle C_{n} S_{n}\right\rangle\right|^{2} /\left.\left(\left.\langle | C_{n}\right|^{2}\right\rangle\langle | S_{n}\right|^{2}\right\rangle$. Two peakes of the population distribution, corresponding to each of the levels, can be seen.

From the analysis of the two-level-band system we have learned that the correlation of the matrix elements of interaction may have an important influence on the population dynamics. These correlations may result in a considerable increase in the decay time as was the case for (4.129), which however remains finite and roughly amounts to the dispersion of the matrix elements $V^{2}\left(1-\left|\left\langle e^{i \chi_{n}-i \theta_{n}}\right\rangle\right|^{2}\right)$ multiplied by the band state density. If the dispersion of the matrix element distribution is of the order of its mean square, the typical decay time does not experience drastic variations, although the distribution of the population over the band can be essentially non-Lorentzian.

### 4.5.2 Two Bands With a Correlated Coupling

Let us now consider another example of a system with correlated matrix elements, namely a system of two bands (4.5) with a Hamiltonian (4.6) where the coupling matrices have the factorized form $V_{n}^{m}=C_{n} \mathcal{V} C_{m}^{*}$ of (4.107). For the sake of simplicity we assume a specific form

$$
\begin{align*}
C_{n} & =\sqrt{\frac{N}{\pi\left(N^{2}+n^{2}\right)}} \\
C_{m} & =\sqrt{\frac{M}{\pi\left(M^{2}+m^{2}\right)}} \tag{4.138}
\end{align*}
$$

of the coefficients entering (4.107) that allows one to perform the consideration analytically.

We note that
a)

b)


Fig. 4.24. (a) Two bands coupled by a factorized interaction. The $N$ levels $|n\rangle$ of the first band located in a strip of width $\sim N \delta_{0}$ interact with $M$ states $|m\rangle$ of the second band within a strip of width $\sim M \delta_{1}$ as if the interaction occurs via a single intermediate state $\left|V_{s}\right\rangle=C_{n}|n\rangle=C_{m}|m\rangle$, which however does not correspond to an energy level. The mean spacing among the neighboring levels of the first and the second bands are $\delta_{0}$ and $\delta_{1}$ respectively. (b) The same problem in another representation, where the states $|b 1\rangle=\sum C_{n}|n\rangle$ and $|b 2\rangle=\sum C_{m}|m\rangle$ are included in the basis set together with the orthogonal linear combinations $\left|n^{\prime}\right\rangle=\sum C_{n}^{n^{\prime}}|n\rangle$ of the states of the first band and $\left|m^{\prime}\right\rangle=\sum C_{m}^{m^{\prime}}|m\rangle$ of the second band. States $|b 1\rangle$ and $|b 2\rangle$ are coupled by the cooperative interaction $\mathcal{V}=V \sqrt{M N}$.

$$
\begin{equation*}
\sum_{n=-\infty}^{\infty}\left|C_{n}\right|^{2}=\sum_{m=-\infty}^{\infty}\left|C_{m}\right|^{2} \simeq 1 \tag{4.139}
\end{equation*}
$$

for large $N$ and $M$. This implies that the rank of the perturbation matrix equals one, and that the matrix itself is proportional to a direct product of two normalized vectors $\left|c_{0}\right\rangle\left\langle c_{1}\right|$ that belong to different bands. The only nonzero eigenvalue $\mathcal{V}$ of $V_{n}^{m}$ corresponds to the cooperative matrix element for the transition between the bands, which couples $\left|c_{0}\right\rangle$ with $\left|c_{1}\right\rangle$.

What is the physical meaning of (4.138)? It tells us that a number of the order of $N$ levels of the first band interact strongly with a number of the order of $M$ levels of the second band, and that the matrix element of the interaction changes smoothly with the state numbers $n$ and $m$. Let us denote by $\delta_{0}$ and $\delta_{1}$ the mean energy spacing between the neighboring levels of the first and the second bands respectively. Then one can say that the levels located in the energy strip $N \delta_{0}$ in the vicinity of $\Delta_{n}=0$ interact strongly with the levels of the second band located in the energy strip $M \delta_{1}$ around the position $\Delta_{m}=0$ of the second band, as shown in Fig. 4.24. One can say that the interaction of the bands occurs via an intermediate state $\left|V_{s}\right\rangle$, which does not correspond to an energy level, being just the only non-zero eigenvector of the interaction matrix (4.106). Projection of this vector onto the Hilbert space of the first band gives the state vector $\left|c_{0}\right\rangle=\sum_{n} C_{n}|n\rangle$, whereas for the Hilbert space of the second band this projection reads $\left|c_{1}\right\rangle=\sum_{m} C_{m}|m\rangle$. These states are responsible for the interaction of the bands. When this interaction is a dipole coupling of atomic multiplets mediated by an external electromagnetic field they are called "bright states", in contrast to the "dark states" that are orthogonal to the "bright" ones and do not interact directly with the field.

## Formal Solution

The Shcrödinger equation (4.5) for the system yields

$$
\begin{align*}
& \psi_{n}(\varepsilon)=\frac{C_{n} \mathcal{V}}{\left(\varepsilon-\Delta_{n}\right)} \sum_{m=-\infty}^{\infty} C_{m} \psi_{m}(\varepsilon)+\frac{i \delta_{n}^{0}}{\left(\varepsilon-\Delta_{n}\right)} \\
& \psi_{m}(\varepsilon)=\frac{C_{m} \mathcal{V}}{\left(\varepsilon-\Delta_{m}\right)} \sum_{n=-\infty}^{\infty} C_{n} \psi_{n}(\varepsilon) \tag{4.140}
\end{align*}
$$

We multiply the first equation by $C_{n}$, the second equation by $C_{m}$, sum up the terms, and obtain the equations

$$
\begin{align*}
& \Phi_{0}(\varepsilon)=\sum_{n=-\infty}^{\infty} \frac{\left|C_{n}\right|^{2} \mathcal{V}}{\left(\varepsilon-\Delta_{n}\right)} \Phi_{1}(\varepsilon)+\frac{i C_{0}}{\left(\varepsilon-\Delta_{0}\right)} \\
& \Phi_{1}(\varepsilon)=\sum_{m=-\infty}^{\infty} \frac{\left|C_{m}\right|^{2} \mathcal{V}}{\left(\varepsilon-\Delta_{m}\right)} \Phi_{0}(\varepsilon), \tag{4.141}
\end{align*}
$$

for the combinations

$$
\begin{align*}
& \Phi_{0}(\varepsilon)=\sum_{n=-\infty}^{\infty} C_{n} \psi_{n}(\varepsilon) \\
& \Phi_{1}(\varepsilon)=\sum_{m=-\infty}^{\infty} C_{m} \psi_{m}(\varepsilon) . \tag{4.142}
\end{align*}
$$

We now substitute (4.138) into the sums entering (4.141), substitute $\Delta_{n}=$ $n \delta_{0}, \Delta_{m}=m \delta_{1}$ and replace the summation by integration in the regime $t<\left(1 / \delta_{0} ; 1 / \delta_{1}\right)$. We arrive at

$$
\begin{align*}
\sum_{n=-\infty}^{\infty} \frac{\left|C_{n}\right|^{2}}{\left(\varepsilon-\Delta_{n}\right)} & =\frac{1}{\pi} \int_{-\infty}^{\infty} d n \frac{N}{N^{2}+n^{2}} \frac{1}{\varepsilon-n \delta_{0}}=\frac{1}{\varepsilon+i N \delta_{0}} \\
\sum_{m=-\infty}^{\infty} \frac{\left|C_{m}\right|^{2}}{\left(\varepsilon-\Delta_{m}\right)} & =\frac{1}{\pi} \int_{-\infty}^{\infty} d m \frac{M}{M^{2}+m^{2}} \frac{1}{\varepsilon-m \delta_{1}}=\frac{1}{\varepsilon+i M \delta_{1}} \tag{4.143}
\end{align*}
$$

Note the physical meaning of these sums. They represent the Fourier transformed autocorrelation, that is the probability amplitude to find the system at the states $\left|c_{0}\right\rangle$ and $\left|c_{1}\right\rangle$, respectively, after a time interval $t$, provided the system had been in these states at $t=0$. Substituting (4.143) into (4.141) we obtain

$$
\begin{align*}
\Phi_{0}(\varepsilon) & =\frac{\mathcal{V}}{\varepsilon+i N \delta_{0}} \Phi_{1}(\varepsilon)+\frac{i}{\varepsilon \sqrt{\pi N}} \\
\Phi_{1}(\varepsilon) & =\frac{\mathcal{V}}{\varepsilon+i M \delta_{1}} \Phi_{0}(\varepsilon) \tag{4.144}
\end{align*}
$$

and solving this system of equations we find

$$
\begin{align*}
& \Phi_{0}(\varepsilon)=\frac{i}{\varepsilon \sqrt{\pi N}}\left(1-\frac{\mathcal{V}}{\varepsilon+i N \delta_{0}} \frac{\mathcal{V}}{\varepsilon+i M \delta_{1}}\right)^{-1} \\
& \Phi_{1}(\varepsilon)=\frac{\mathcal{V}}{\varepsilon+i M \delta_{1}} \frac{i}{\varepsilon \sqrt{\pi N}}\left(1-\frac{\mathcal{V}}{\varepsilon+i N \delta_{0}} \frac{\mathcal{V}}{\varepsilon+i M \delta_{1}}\right)^{-1} \tag{4.145}
\end{align*}
$$

which yields, with the allowance of (4.140),

$$
\begin{align*}
\psi_{n}(\varepsilon)= & \frac{i \delta_{n}^{0}}{\varepsilon}+\frac{C_{n} \mathcal{V}}{\varepsilon-\Delta_{n}} \frac{\mathcal{V}}{\varepsilon+i M \delta_{1}} \frac{i}{\varepsilon \sqrt{\pi N}} \\
& \times\left(1-\frac{\mathcal{V}}{\varepsilon+i N \delta_{0}} \frac{\mathcal{V}}{\varepsilon+i M \delta_{1}}\right)^{-1} \\
\psi_{m}(\varepsilon)= & \frac{C_{m} \mathcal{V}}{\varepsilon-\Delta_{m}} \frac{i}{\varepsilon \sqrt{\pi N}}\left(1-\frac{\mathcal{V}}{\varepsilon+i N \delta_{0}} \frac{\mathcal{V}}{\varepsilon+i M \delta_{1}}\right)^{-1} . \tag{4.146}
\end{align*}
$$

## Distribution of the Population over the Band

Let us concentrate on the last equation (4.146) and consider the regime $\left(1 / N \delta_{0} ; 1 / M \delta_{1}\right) \ll t \ll\left(1 / \delta_{0} ; 1 / \delta_{1}\right)$ when only the poles at the real points $\varepsilon=\Delta_{m} ; \varepsilon=0$ are important in the inverse Fourier transformation, whereas two other poles at the imaginary points $\varepsilon=i N \delta_{0}$ and $\varepsilon=i M \delta_{1}$ yield in this regime, exponentially vanishing contributions.

The inverse Fourier transformation gives the time dependent amplitude $\psi_{m}(t)$

$$
\begin{align*}
\psi_{m}(t)= & \frac{-C_{m} \mathcal{V}}{\Delta_{m} \sqrt{\pi N}}\left[\frac{N \delta_{0} M \delta_{1}}{N \delta_{0} M \delta_{1}+\mathcal{V}^{2}}\right. \\
& \left.-\frac{e^{i \Delta_{m} t}\left(\Delta_{m}-i N \delta_{0}\right)\left(\Delta_{m}-i M \delta_{1}\right)}{\left(\Delta_{m}-i N \delta_{0}\right)\left(\Delta_{m}-i M \delta_{1}\right)-\mathcal{V}^{2}}\right] \tag{4.147}
\end{align*}
$$

and the population $\rho_{m}(t)=\left|\psi_{m}(t)\right|^{2}$ reads

$$
\begin{align*}
\rho_{m}(t)= & \frac{1 / N M}{1+y^{2} \Delta^{2}} \frac{w}{\Delta^{2} \pi^{2}}\left|\frac{e^{i \Delta \tau}(\Delta-i y)(\Delta-i / y)}{(\Delta-i y)(\Delta-i / y)-w}-\frac{1}{1+w}\right|^{2} \\
= & \frac{1 / N M}{1+y^{2} \Delta^{2}} \frac{w}{\pi^{2}} \left\lvert\, \frac{e^{i \Delta \tau}-1}{\Delta} \frac{(\Delta-i y)(\Delta-i / y)}{(\Delta-i y)(\Delta-i / y)-w}\right. \\
& +\left.\frac{w}{(1+w)} \frac{\Delta-(i y+i / y)}{((\Delta-i y)(\Delta-i / y)-w)}\right|^{2} \tag{4.148}
\end{align*}
$$

where $w=\mathcal{V}^{2} / N \delta_{0} M \delta_{1} ; \Delta=\Delta_{m} / \sqrt{N \delta_{0} M \delta_{1}} ; y=\sqrt{N \delta_{0} / M \delta_{1}} ; \tau=$ $t \sqrt{N} \delta_{0} M \delta_{1}$ and we have employed (4.138). One sees the important peculiarity of this system - no stationary distribution over the levels of the second
band arises after the relaxation time $t \sim\left(1 / N \delta_{0} ; 1 / M \delta_{1}\right)$, since the righthand side of (4.147) remains a function of time.

In the long-time limit, the right-hand side of (4.148) can be considered as consisting of two parts. The well-localized first part arising in the vicinity of the point $\Delta=0$ from the terms containing the combination $\left(e^{i \Delta \tau}-1\right) / \Delta$ and its square is given by


Fig. 4.25. Two bands coupled by a factorized interaction. (a) Distribution of population over the second band for $y=1.2$ for three different times $\tau=1,3,9$ in the regime $\operatorname{Im} \sqrt{r} \neq 0$. One sees that the time-dependent part of the population becomes larger. (b) The same as (a) for $\operatorname{Im} \sqrt{r}=0$. (c) The time-independent part of the population distribution as a function of $w$. One sees the cross-over from the regime $\operatorname{Im} \sqrt{r}=0$ where the distribution is localized around $\Delta=0$, to the regime $\operatorname{Im} \sqrt{r} \neq 0$ where tree peaks of the distribution are present. (d) The time-dependent part of the distribution becomes more and more narrow with the course of time and contains a larger and larger part of the population. The total contribution of this spike increases as $\sim \tau$.

$$
\begin{align*}
\rho_{m}^{(1)}(t) \simeq & \frac{1}{N M} \frac{2 w^{2}\left(y+y^{-1}\right)}{\pi^{2}(1+w)^{3}} \frac{\sin (t \Delta)}{\Delta} \\
& -\frac{1}{N M} \frac{2 w}{\pi^{2}(1+w)^{2}} \frac{(\cos (t \Delta)-1)}{\Delta^{2}} \tag{4.149}
\end{align*}
$$

where as $t \rightarrow \infty$ one recognizes the Dirac $\delta$-function $\delta(\Delta)$ which originates from the first term on the right-hand side of (4.149) linear in $\left(e^{i \Delta \tau}-1\right) / \Delta$
and the dominating contribution, $\sim \delta(\Delta) t$, resulting from the second term quadratic in $\left(e^{i \Delta \tau}-1\right) / \Delta$. The time-independent second part

$$
\begin{equation*}
\rho_{m}^{(2)}(t) \simeq \frac{1 / N M \pi^{2}}{1+y^{2} \Delta^{2}} \frac{w^{3}}{(1+w)}\left|\frac{\Delta-(i y+i / y)}{((\Delta-i y)(\Delta-i / y)-w)}\right|^{2} \tag{4.150}
\end{equation*}
$$

emerges from (4.148) when at $\Delta \neq 0$ and $\Delta t \gg 1$ we omit the rapidly oscillating term $\left(e^{i \Delta \tau}-1\right) / \Delta$.

The time-independent part (4.150) of the population distribution can be represented as a sum of three Lorentzian profiles - one peak of width $1 /|y|$ located at $\Delta=0$, and two other peaks, of width $a=\left|\operatorname{Re}_{1}\right|$ located at the points $\Delta= \pm b= \pm \operatorname{Im} r_{1}$, where

$$
\begin{equation*}
r_{1,2}=a \pm i b=\frac{y+y^{-1}}{2} \pm \sqrt{\left(\frac{y-y^{-1}}{2}\right)^{2}-w} \tag{4.151}
\end{equation*}
$$

is the solution of the bi-quadratic equation $|(i x-i y)(i x-i / y)-w|^{2}=0$ with respect to $\Delta=i x$, corresponding to the poles of (4.150). For real and positive $r$ all three of these peaks are located at the point $\Delta=0$, but have nevertheless different widths. In Fig. 4.25 we show the population distribution profiles as a function of the detuning for different regimes.

## Total Populations of the Bands

Multiplying the probability amplitude $\psi_{m}(\varepsilon)$ of (4.146) by its complex conjugate one finds the total population of the second band by integrating over $d m=d \Delta / \delta_{1}$. After substitution of $C_{m}$ from (4.138) and employing the notation of (4.148)

$$
\begin{align*}
\rho_{1}(\varepsilon, \xi)= & \int \frac{d \Delta}{\pi^{2} N} \frac{1 / y}{\left(1 / y^{2}+\Delta^{2}\right)} \frac{w}{(\varepsilon-\Delta)(\xi-\Delta)} \frac{1}{\varepsilon \xi} \\
& \left(1+\frac{w}{\left(\varepsilon-i r_{1}\right)\left(\varepsilon-i r_{2}\right)}\right)\left(1+\frac{w}{\left(\xi+i r_{1}\right)\left(\xi+i r_{2}\right)}\right) \tag{4.152}
\end{align*}
$$

one obtains

$$
\begin{align*}
\rho_{1}(\varepsilon, \xi)= & \frac{w}{N \pi}\left(1+\frac{2 i}{y(\varepsilon-\xi)}\right) \\
& \frac{1}{\varepsilon\left(\varepsilon+\frac{i}{y}\right)}\left(1+\frac{w}{(\varepsilon-b+a i)(\varepsilon+b+a i)}\right) \\
& \frac{1}{\xi\left(\xi-\frac{i}{y}\right)}\left(1+\frac{w}{(\xi-b-a i)(\xi+b-a i)}\right) . \tag{4.153}
\end{align*}
$$

By analogy to (4.75) we note that the factor $i /(\varepsilon-\xi)$ corresponds to the time integration whereas the combinations on the second and the third lines
yield, after the inverse Fourier transformation, two complex conjugate time dependent functions $\phi_{\varepsilon}(t)$ and $\phi_{\xi}(t)$ respectively, which read

$$
\begin{align*}
\phi_{\varepsilon, \xi}(t)= & \int_{C} \frac{e^{-i \varepsilon t}}{\varepsilon\left(\varepsilon+\frac{i}{y}\right)}\left(1+\frac{w}{(\varepsilon-b+a i)(\varepsilon+b+a i)}\right) \\
= & y\left(\frac{w}{b^{2}+a^{2}}-1\right)-\frac{w \sin (t b \pm \varphi) e^{-a t}}{b|b-a i|\left|b-a i+\frac{i}{y}\right|} \\
& +y\left(1-\frac{w}{b^{2}-(a-1 / y)^{2}}\right) e^{-t / y} \tag{4.154}
\end{align*}
$$

where $\varphi=\arg (b-a i)+\arg \left(b-a i+\frac{i}{y}\right)$. Therefore the time-dependent population $\rho_{1}(t)$ of the second band can be given in the form

$$
\begin{equation*}
\rho_{1}(\tau)=\frac{w}{N \pi}\left(\left|\phi_{\varepsilon}(\tau)\right|^{2}+\frac{2}{y} \int_{0}^{\tau}\left|\phi_{\varepsilon}(t)\right|^{2} d t\right) \tag{4.155}
\end{equation*}
$$

which after straightforward calculations yields

$$
\begin{align*}
\rho_{1}(\tau)= & \frac{2 w}{N(1+w) \pi^{2}}\left[\frac{2 \tau y}{w}+\frac{w^{2}+w\left(2+y^{2}\right)^{2}-2\left(y^{2}+y^{4}\right)}{w(1+w)\left(1+y^{2}\right)}\right. \\
& -\frac{2 y^{2}(1+w)\left(y^{2}-1\right) e^{-2 a \tau}}{\left(1+y^{2}\right)\left(1-2(1+2 w) y^{2}+y^{4}\right)} \\
& +\frac{2\left(y^{2}-w\left(2+y^{2}\right)\right)}{(1+w)} e^{-a \tau} \cos (b \tau) \\
& +\frac{2 w^{2} y^{2}+y^{2}\left(y^{2}-1\right)-w\left(2+3 y^{2}+y^{4}\right)}{b(1+w) y w} e^{-a \tau} \sin (b \tau) \\
& +\frac{2 w^{2} y^{2}+y^{4}\left(y^{2}-1\right)-w\left(1+y^{2}+4 y^{4}\right)}{(1+w)\left(1-(2+4 w) y^{2}+y^{4}\right)} e^{-2 a \tau} \cos (2 b \tau) \\
& \left.-\frac{w+2 w y^{2}-y^{4}}{2 b(1+w) y} e^{-2 a \tau} \sin (2 b \tau)\right] . \tag{4.156}
\end{align*}
$$

One sees that the total population demonstrates an exponential decay along with dying oscillations. We should note, however, that for the initial conditions chosen, these terms play a minor role, remaining of the order of inverse numbers $1 / N ; 1 / M$, of levels participating in the population dynamics. The dominating role is played by the linear term

$$
\begin{equation*}
\rho_{1}(\tau) \simeq 4 \tau y / N(1+w) \pi^{2}=4 t \delta_{0} /\left(1+\mathcal{V}^{2} / N \delta_{0} M \delta_{1}\right) \pi^{2} \tag{4.157}
\end{equation*}
$$

which however also remains small as compared to unity within the domain $t \ll g \sim 1 / \delta_{0}$ of applicability of the continuous-band model.

## Representation by Two Level-Band Systems

The behavior of the two-band system with factorized interaction matrix elements is revealed when one considers the matrix structure of the corresponding Hamiltonian. In Fig. 4.24 we illustrate the principal features of such a structure, by representing the two-band system as a combination of two levelband systems interacting via coupled levels. Indeed, the Hamiltonian of a two band system with the factorized interaction $\hat{V}=\left|C_{n}\right\rangle \mathcal{V}\left\langle C_{m}\right|$

$$
\hat{H}=\left(\begin{array}{cc}
\sum_{m}|m\rangle \Delta_{m}\langle m| & \left|c_{1}\right\rangle \mathcal{V}\left\langle c_{0}\right|  \tag{4.158}\\
& \left|c_{0}\right\rangle \mathcal{V}\left\langle c_{1}\right| \\
\sum_{n}|n\rangle \Delta_{n}\langle n|
\end{array}\right)=\left(\begin{array}{cc}
\Delta_{m} \delta_{m}^{m^{\prime}} & C_{n}^{*} \mathcal{V} C_{m} \\
C_{m}^{*} \mathcal{V} C_{n} & \Delta_{n} \delta_{n}^{n^{\prime}}
\end{array}\right)
$$

can be rewritten in another representation, which comprises the states $\left|c_{0}\right\rangle$ and $\left|c_{1}\right\rangle$ as the basis states together with the other orthogonal basis states $\{\nu\}$ and $\{\mu\}$, with $\mu, \nu=1,2, \ldots$ The Hamiltonians of the first $\sum_{n}|n\rangle \Delta_{n}\langle n|$ and the second $\sum_{m}|m\rangle \Delta_{m}\langle m|$ bands are not diagonal in such a representation, but via an appropriate choice of the sets $|\nu\rangle$ and $|\mu\rangle$ one can diagonalize the blocks of the Hamiltonian corresponding to these states. This yields the representation

$$
\hat{H}=\left(\begin{array}{cccc}
|\mu\rangle \Delta_{\mu}^{\prime}\langle\mu| & |\mu\rangle H_{\mu}\left\langle c_{1}\right| & 0 & 0  \tag{4.159}\\
\left|c_{1}\right\rangle H_{\mu}^{*}\langle\mu| & \left|c_{1}\right\rangle E_{1}\left\langle c_{1}\right| & \left|c_{1}\right\rangle \mathcal{V}\left\langle c_{0}\right| & 0 \\
0 & \left|c_{0}\right\rangle \mathcal{V}\left\langle c_{1}\right| & \left|c_{0}\right\rangle E_{0}\left\langle c_{0}\right| & |\nu\rangle H_{\nu}\left\langle c_{0}\right| \\
0 & 0 & \left|c_{0}\right\rangle H_{\nu}^{*}\left\langle n^{\prime}\right| & |\nu\rangle \Delta_{\nu}^{\prime}\langle\nu|
\end{array}\right),
$$

in Dirac notation, where summation over $\nu$ and $\mu$ is implicit, or explicitly in the matrix form

$$
\hat{H}=\left(\begin{array}{cccccccccc}
. & . & . & . & . & . & . & . & . & .  \tag{4.160}\\
. & \Delta_{\mu-1}^{\prime} & 0 & 0 & H_{\mu-1} & 0 & 0 & 0 & 0 & . \\
. & 0 & \Delta_{\mu}^{\prime} & 0 & H_{\mu} & 0 & 0 & 0 & 0 & . \\
. & 0 & 0 & \Delta_{\mu+1}^{\prime} & H_{\mu+1} & 0 & 0 & 0 & 0 & . \\
. H_{\mu-1}^{*} & H_{\mu}^{*} & H_{\mu+1}^{*} & E_{1} & \mathcal{V} & 0 & 0 & 0 & . \\
. & 0 & 0 & 0 & \mathcal{V} & E_{0} & H_{\nu+1} & H_{\nu} & H_{\nu-1} & . \\
. & 0 & 0 & 0 & 0 & H_{\nu+1}^{*} & \Delta_{\nu+1}^{\prime} & 0 & 0 & . \\
. & 0 & 0 & 0 & 0 & H_{\nu}^{*} & 0 & \Delta_{\nu}^{\prime} & 0 & . \\
. & 0 & 0 & 0 & 0 & H_{\nu-1}^{*} & 0 & 0 & \Delta_{\nu-1}^{\prime} & . \\
. & . & . & . & . & . & . & . & . & .
\end{array}\right),
$$

where $E_{0}=\sum_{n}\left|C_{n}\right|^{2} \Delta_{n}$ and $E_{1}=\sum_{m}\left|C_{m}\right|^{2} \Delta_{m}$.
The precise positions of the energies $\Delta_{\nu}^{\prime}$ and $\Delta_{\mu}^{\prime}$ are not important in the regime $t \ll g$, and can, as earlier, be taken as equidistant with the spacing among the neighboring levels $\delta_{0}$ and $\delta_{1}$ respectively, while the new non-diagonal matrix elements $H_{\mu}$ and $H_{\nu}$ can be estimated with the help of the consideration on p. 75 of the number of levels in resonance: We express
the number $N$ of levels of the first band accessible from the state $\left|c_{0}\right\rangle$ as $\left.\left.\langle | H_{\nu}\right|^{2}\right\rangle / \delta_{0}^{2}=N$ and by analogy $\left.\left.\langle | H_{\mu}\right|^{2}\right\rangle / \delta_{1}^{2}=M$ for the second band, which yields $H_{\nu} \sim \delta_{0} \sqrt{N}$ and $H_{\mu} \sim \delta_{1} \sqrt{M}$ respectively.

We note that the combinations $\Phi_{0}(\varepsilon)$ and $\Phi_{1}(\varepsilon)$ of (4.142) now have a clear physical meaning, representing the probability amplitudes to be in the states $\left|c_{0}\right\rangle$ and $\left|c_{1}\right\rangle$, whereas (4.144) is the Fourier form of the Schrödinger equation

$$
\begin{align*}
& \varepsilon \Phi_{0}(\varepsilon)=i N \delta_{0} \Phi_{0}(\varepsilon)+\mathcal{V} \Phi_{1}(\varepsilon)+\frac{i\left(\varepsilon+i N \delta_{0}\right)}{\varepsilon \sqrt{\pi N}} \\
& \varepsilon \Phi_{1}(\varepsilon)=\mathcal{V} \Phi_{0}(\varepsilon)+i M \delta_{1} \Phi_{1}(\varepsilon) \tag{4.161}
\end{align*}
$$

for these amplitudes allowing for the irreversible decay with the rates $N \delta_{0}$ and $M \delta_{1}$ to the continua of the first and the second bands respectively. In the inhomogeneous term on the right-hand side of (4.161) one also recognizes (see p. 83) a small initial population amplitude $1 / \sqrt{\pi N}$ of the state $\left|c_{0}\right\rangle$ at $t=0$, and a steady state amplitude injection $\Pi(t)=\Theta(t) \sqrt{N} \delta_{0} / \sqrt{\pi}$ corresponding to the term $\sqrt{N} \delta_{0} / \varepsilon \sqrt{\pi}$.

The origin of the terms exponentially decaying, oscillating, and linearly increasing with time in (4.156) for the population is now clear: the linearly increasing term originates from the steady state population injection, which is becoming localized with the course of time near the energy $\Delta=0$, whereas the terms exponentially decreasing in time result from the irreversible (on the scale $t \ll g$ ) decay of the small initial population of the state $\left|c_{0}\right\rangle$. The decay occurs either with or without oscillations depending on the relation between the coupling strength $\mathcal{V}$ and the typical damping rates $N \delta_{0}$ and $M \delta_{1}$ similar to the dynamics of a damped pendulum. Large coupling corresponds to the oscillatory regime, which results in satellite spikes in the population distribution over the band at the positions $\Delta \sim \pm \mathcal{V}$, seen in Fig. 4.25(a,b,c).

But where does the initial population of the state $\left|c_{0}\right\rangle$ and the population amplitude supply come from? To answer this question, one has to note that the initial condition $\psi_{n}(t=0)=\delta_{n}^{0}, \psi_{m}(t=0)=0$ in the representation of the basis states $|\nu\rangle$ and $|\mu\rangle$ is determined by the amplitudes $\psi_{\nu}(t=0)$ and $\psi_{\mu}(t=0)$ and reads $\psi_{\nu}(t=0)=\langle\nu \mid n=0\rangle, \psi_{\mu}(t=0)=0$. This yields a small initial amplitude $\left\langle c_{0} \mid n=0\right\rangle=C_{0}=1 / \sqrt{\pi N}$ for the state $\left|c_{0}\right\rangle$, whereas the main part $(\sim 1)$ of the initial population amplitude remains in the band states $|\nu\rangle$ either precisely at the state $|\nu=0\rangle$ or at a small group of states located in an extremely narrow strip around it. Being coupled to the shortlived state $\left|c_{0}\right\rangle$ by a relatively small matrix element $H_{\nu=0} \sim \delta_{0} \sqrt{N}$, the state $|\nu=0\rangle$ supplies a small population flux $\Pi \sim\left|H_{\nu=0}\right|^{2} /\left(N \delta_{0}+M \delta_{1}\right)$ given by the ratio of the squared coupling to the decay rate of $\left|c_{0}\right\rangle$. (Note that this ratio represents the cooperative matrix element of coupling squared divided by the width of the band discussed on p. 76.) For $N \delta_{0} \sim M \delta_{1}$ this yields the estimate $\Pi \sim \rho_{1}(t) / t \sim \delta_{0}$ which is consistent with the result of (4.157).

One sees that the correlations of the interaction matrix elements play an important role and may have a significant effect on the population dynamics in a multilevel quantum system. The Fano problem discussed in Sect.3.2.2 and the problems that have been considered in this section show that the correlations manifest themselves in non-Lorentzian profiles of the population distribution over the band and in the non-exponential time dependences of the total populations of the bands.

### 4.5.3 Regime of Stabilization for the Correlated Coupling

An interesting regime exists in a system consisting of an equidistant discrete band coupled to a continuum by a correlated coupling $V$. It turns out that for time $t$ considerably exceeding the inverse spacing $1 / \delta$ between the neighboring discrete levels, the rate of the band population decay to continuum decreases with the increase of $V$. We consider the continuum with a state density $g=1 / \delta_{1}$ as a band of $N$ discrete levels in the limit $N \rightarrow \infty, \delta_{1} \rightarrow 0$ and employ (4.140) with $C_{n}=1 / \sqrt{N}$ and $C_{m}=1 / \sqrt{M}$, where $M$ is the number of levels in the disctete band. This yields

$$
\begin{align*}
& \psi_{n}(\varepsilon)=\frac{\mathcal{V}}{\sqrt{M N( }\left(\varepsilon-\Delta_{n}\right)} \sum_{m=-\infty}^{\infty} \psi_{m}(\varepsilon)+\frac{i \Pi(\varepsilon)}{\sqrt{N}\left(\varepsilon-\Delta_{n}\right)}, \\
& \psi_{m}(\varepsilon)=\frac{\mathcal{V}}{\sqrt{M N}\left(\varepsilon-\Delta_{m}\right)} \sum_{n=-\infty}^{\infty} \psi_{n}(\varepsilon), \tag{4.162}
\end{align*}
$$

where $\Pi(\varepsilon)=\exp \left(-\varepsilon^{2} \tau^{2} / 4\right)$ is the Fourier transformed external injection (3.58) to the continuum and $\mathcal{V}=V \sqrt{M N}$.

By analogy to (4.140)-(4.146) with (4.143) of the form

$$
\begin{align*}
\int_{-\infty}^{\infty} d n \frac{1}{\varepsilon-n \delta_{1}} & =i \pi g \\
\sum_{m=-\infty}^{\infty} \frac{1}{\varepsilon-m \delta} & =\frac{\pi}{\delta} \cot \frac{\varepsilon \pi}{\delta} \tag{4.163}
\end{align*}
$$

suggested by (3.76), we arrive at

$$
\begin{equation*}
\psi_{m}(\varepsilon)=\frac{-\pi g \Pi(\varepsilon) V}{\left(\varepsilon-\Delta_{m}\right)}\left(1-\frac{i \pi^{2} g V^{2}}{\delta} \cot \frac{\varepsilon \pi}{\delta}\right)^{-1} . \tag{4.164}
\end{equation*}
$$

We determine the level population

$$
\begin{align*}
\rho_{m}(\varepsilon, \xi)= & (\pi g V)^{2} \frac{\Pi(\varepsilon) \Pi(\xi)}{\left(\varepsilon-\Delta_{m}\right)\left(\xi-\Delta_{m}\right)} \\
& \left(1-\frac{i \pi^{2} g V^{2}}{\delta} \cot \frac{\varepsilon \pi}{\delta}\right)^{-1}\left(1+\frac{i \pi^{2} g V^{2}}{\delta} \cot \frac{\xi \pi}{\delta}\right)^{-1} \tag{4.165}
\end{align*}
$$

and find the total population $\rho_{1}(\varepsilon, \xi)=\sum_{m} \rho_{m}(\varepsilon, \xi)$ of the discrete band

$$
\begin{align*}
\rho_{1}(\varepsilon, \xi)= & -(\pi g V)^{2} \Pi(\varepsilon) \Pi(\xi) \frac{\cot \pi \xi / \delta-\cot \pi \varepsilon / \delta}{(\xi-\varepsilon)} \\
& \frac{\pi}{\delta}\left(1-\frac{i \pi^{2} g V^{2}}{\delta} \cot \frac{\varepsilon \pi}{\delta}\right)^{-1}\left(1+\frac{i \pi^{2} g V^{2}}{\delta} \cot \frac{\xi \pi}{\delta}\right)^{-1} \tag{4.166}
\end{align*}
$$

For large $g V^{2} / \delta$ and for long times $t \gg \pi / \delta$ the main contribution to the inverse Fourier transformation comes from the vicinities of the points where $\xi \simeq \varepsilon \simeq \pi / 2+n \pi$, and after having performed the replacements $\xi \rightarrow \xi+$ $\pi / 2+n \pi$ and $\varepsilon \rightarrow \varepsilon+\pi / 2+n \pi$ we arrive at

$$
\begin{align*}
\rho_{1}(\varepsilon, \xi)= & -(\pi g V)^{2} \sum_{n} \Pi\left(n \delta_{1}\right) \Pi\left(n \delta_{1}\right) \frac{\pi}{\delta_{1}} \frac{\pi}{\delta_{1}} \\
& \left(1-\frac{i \pi^{3} g V^{2}}{\delta^{2}} \varepsilon\right)^{-1}\left(1+\frac{i \pi^{3} g V^{2}}{\delta^{2}} \xi\right)^{-1} \tag{4.167}
\end{align*}
$$

and hence

$$
\begin{equation*}
\rho_{1}(t)=\frac{4 \delta^{2}}{V^{2}} e^{-2 t \delta^{2} / \pi^{3} g V^{2}} \sum_{n} e^{-(n \delta \tau)^{2} / 2}=\frac{4 g \delta \sqrt{2 \pi}}{\tau W} e^{-2 t \delta^{2} / \pi^{2} W} \tag{4.168}
\end{equation*}
$$

One sees that the stronger the regular decay rate $2 W=2 \pi g V^{2}$ of the band levels to the continuum at $t<1 / \delta$, the slower is this decay rate $2 \delta^{2} / \pi^{2} W$ in the regime $t \gg \pi / \delta$ under consideration. This phenomenon has been observed in Rydberg atoms interacting with a strong laser field and is known as strong field stabilization of atoms.

### 4.5.4 Correlation Between the Mean Squared Coupling and the Energy Position

Let us now consider a different example of the correlation, when the interaction matrix elements $V_{n}^{m}$ do not correlate among themselves, that is $\left\langle V_{n}^{m}\right\rangle=0$, but their mean square values $\left.v_{n}^{m}=\left.\langle | V_{n}^{m}\right|^{2}\right\rangle$ are different for different transitions, and therefore they correlate with the energy positions $\Delta_{m}$ and $\Delta_{n}$ of both the coupled states $|m\rangle$ and $|n\rangle$. All the results of Sect. 4.1 remain valid for this case if one does not put the mean squared interactions in front of the summations symbols, such that (4.20) and (4.21) read

$$
\begin{align*}
X_{n}(\varepsilon) & =\frac{1}{\varepsilon-\Delta_{n}-\sum_{m} v_{n}^{m} X_{m}(\varepsilon)} \\
X_{m}(\varepsilon) & =\frac{1}{\varepsilon-\Delta_{m}-\sum_{n} v_{m}^{n} X_{n}(\varepsilon)} \tag{4.169}
\end{align*}
$$

while (4.26) takes the form

$$
\begin{align*}
\rho_{m}(\varepsilon, \xi)= & \frac{1}{\varepsilon-\Delta_{m}-\sum_{n} v_{m}^{n} X_{n}(\varepsilon)} \frac{1}{\xi-\Delta_{m}-\sum_{n} v_{m}^{n} X_{n}(\xi)} \sum_{n} v_{m}^{n} \rho_{n}(\varepsilon, \xi) \\
\rho_{n}(\varepsilon, \xi)= & \frac{1}{\varepsilon-\Delta_{n}-\sum_{m} v_{n}^{m} X_{m}(\varepsilon)} \frac{1}{\xi-\Delta_{n}-\sum_{m} v_{n}^{m} X_{m}(\xi)} \sum_{m} v_{n}^{m} \rho_{m}(\varepsilon, \xi) \\
& +\delta_{0}^{n} \frac{1}{\xi-\sum_{m} v_{0}^{m} X_{m}(\xi)} \frac{1}{\varepsilon-\sum_{m} v_{0}^{m} X_{m}(\varepsilon)} . \tag{4.170}
\end{align*}
$$

We concentrate here on the case of a factorized form of the mean squared interaction $v_{n}^{m}=c_{n} V^{2} c_{m}$ and two identical bands for which one obtains the relations

$$
\begin{align*}
X_{n}(\varepsilon) & =X_{m}(\varepsilon)=\frac{1}{\varepsilon-\Delta_{m}-c_{m} V^{2} s(\varepsilon)}  \tag{4.171}\\
s(\varepsilon) & =\sum_{n} c_{n} X_{n}(\varepsilon) \tag{4.172}
\end{align*}
$$

yielding the renormalization equation

$$
\begin{equation*}
s(\varepsilon)=\sum_{n} \frac{c_{n}}{\varepsilon-\Delta_{n}-c_{n} V^{2} s(\varepsilon)}, \tag{4.173}
\end{equation*}
$$

that has much in common with (4.22), although differs from it.
Note that $s(\xi)$ is the complex conjugate function of $s(\varepsilon)$, that is $s(\xi)=$ $\left.s^{*}(\varepsilon)\right|_{\varepsilon=\xi}$. For this case (4.26) reads

$$
\begin{align*}
\rho_{m}(\varepsilon, \xi)= & \frac{1}{\varepsilon-\Delta_{m}-c_{m} V^{2} s(\varepsilon)} \frac{1}{\xi-\Delta_{m}-c_{m} V^{2} s(\xi)} c_{m} V^{2} \widetilde{\rho}_{0}(\varepsilon, \xi) \\
\rho_{n}(\varepsilon, \xi)= & \frac{1}{\varepsilon-\Delta_{n}-c_{n} V^{2} s(\varepsilon)} \frac{1}{\xi-\Delta_{n}-c_{n} V^{2} s(\xi)} c_{n} V^{2} \widetilde{\rho}_{1}(\varepsilon, \xi) \\
& +\delta_{0}^{n} \frac{1}{\xi-c_{0} V^{2} s(\xi)} \frac{1}{\varepsilon-c_{0} V^{2} s(\varepsilon)} \tag{4.174}
\end{align*}
$$

where $\widetilde{\rho}_{0}(\varepsilon, \xi)=\sum_{n} c_{n} \rho_{n}(\varepsilon, \xi)$ and $\widetilde{\rho}_{1}(\varepsilon, \xi)=\sum_{m} c_{m} \rho_{m}(\varepsilon, \xi)$. Multiplication of the first and the second parts of (4.174) by $c_{m}$ and $c_{n}$ respectively followed by summation over all $m$ and $n$ results in the system of equations

$$
\begin{align*}
\widetilde{\rho}_{1}(\varepsilon, \xi) & =\widetilde{F}(\varepsilon, \xi) \widetilde{\rho}_{0}(\varepsilon, \xi) \\
\widetilde{\rho}_{0}(\varepsilon, \xi) & =\widetilde{F}(\varepsilon, \xi) \widetilde{\rho}_{1}(\varepsilon, \xi)+\frac{c_{0}}{\xi-c_{0} V^{2} s(\xi)} \frac{1}{\varepsilon-c_{0} V^{2} s(\varepsilon)} \tag{4.175}
\end{align*}
$$

for the auxiliary quantities $\widetilde{\rho}_{0}(\varepsilon, \xi)$ and $\widetilde{\rho}_{1}(\varepsilon, \xi)$ which has the solution

$$
\begin{align*}
& \widetilde{\rho}_{0}(\varepsilon, \xi)=\frac{-c_{0}}{\left(\widetilde{F}(\varepsilon, \xi)^{2}-1\right)\left(\xi-c_{0} V^{2} s(\xi)\right)\left(\varepsilon-c_{0} V^{2} s(\varepsilon)\right)} \\
& \widetilde{\rho}_{1}(\varepsilon, \xi)=\frac{-\widetilde{F}(\varepsilon, \xi) c_{0}}{\left(\widetilde{F}(\varepsilon, \xi)^{2}-1\right)\left(\xi-c_{0} V^{2} s(\xi)\right)\left(\varepsilon-c_{0} V^{2} s(\varepsilon)\right)} . \tag{4.176}
\end{align*}
$$

After substitution into (4.174) this yields, for the total populations,

$$
\begin{align*}
\rho_{1}(\varepsilon, \xi) & =F(\varepsilon, \xi) \widetilde{\rho}_{0}(\varepsilon, \xi) \\
\rho_{0}(\varepsilon, \xi) & =F(\varepsilon, \xi) \widetilde{\rho}_{1}(\varepsilon, \xi)+\frac{1}{\xi-c_{0} V^{2} s(\xi)} \frac{1}{\varepsilon-c_{0} V^{2} s(\varepsilon)} \tag{4.177}
\end{align*}
$$

Here we have employed the notation

$$
\begin{align*}
& \widetilde{F}(\varepsilon, \xi)=\sum_{m} \frac{V c_{m}}{\varepsilon-\Delta_{m}-c_{m} V^{2} s(\varepsilon)} \frac{V c_{m}}{\xi-\Delta_{m}-c_{m} V^{2} s(\xi)} \\
& F(\varepsilon, \xi)=\sum_{m} \frac{V^{2}}{\varepsilon-\Delta_{m}-c_{m} V^{2} s(\varepsilon)} \frac{c_{m}}{\xi-\Delta_{m}-c_{m} V^{2} s(\xi)} \tag{4.178}
\end{align*}
$$

Note that the sum $\rho_{+}(\varepsilon, \xi)=\rho_{1}(\varepsilon, \xi)+\rho_{0}(\varepsilon, \xi)$ can be written in the form

$$
\begin{align*}
\rho_{+}(\varepsilon, \xi) & =F(\varepsilon, \xi)\left(\widetilde{\rho}_{0}(\varepsilon, \xi)+\widetilde{\rho}_{1}(\varepsilon, \xi)\right)+\frac{1}{\xi-c_{0} V^{2} s(\xi)} \frac{1}{\varepsilon-c_{0} V^{2} s(\varepsilon)} \\
& =\frac{1-\widetilde{F}(\varepsilon, \xi)-c_{0} F(\varepsilon, \xi)}{1-\widetilde{F}(\varepsilon, \xi)} \frac{1}{\xi-c_{0} V^{2} s(\xi)} \frac{1}{\varepsilon-c_{0} V^{2} s(\varepsilon)} \tag{4.179}
\end{align*}
$$

After the inverse Fourier transformation it should be identical to 1, which implies a hidden functional relation among $\widetilde{F}(\varepsilon, \xi), F(\varepsilon, \xi), s(\xi)$ and $s(\varepsilon)$, and indeed, by taking apart the fractional factors in the first part of (4.178) and with allowance for (4.173) one obtains, after straightforward calculations

$$
\begin{equation*}
1-\widetilde{F}(\varepsilon, \xi)=\frac{\varepsilon-\xi}{s(\varepsilon)-s(\xi)} \frac{F(\varepsilon, \xi)}{V^{2}} \tag{4.180}
\end{equation*}
$$

which yields

$$
\begin{align*}
\rho_{+}(\varepsilon, \xi) & =\frac{\varepsilon-\xi-c_{0} V^{2}(s(\varepsilon)-s(\xi))}{\varepsilon-\xi} \frac{1}{\xi-c_{0} V^{2} s(\xi)} \frac{1}{\varepsilon-c_{0} V^{2} s(\varepsilon)} \\
& =\frac{1}{\varepsilon-\xi} \frac{1}{\xi-c_{0} V^{2} s(\xi)}-\frac{1}{\varepsilon-\xi} \frac{1}{\varepsilon-c_{0} V^{2} s(\varepsilon)} \tag{4.181}
\end{align*}
$$

Indeed, consider the first fraction in the last part of the equation. In the inverse Fourier transformation it corresponds to the term, which as a function of $\varepsilon$, has only one pole at the point $\varepsilon=\xi$. We therefore integrate this term over $\varepsilon$ first, thus removing the time dependence. The analytical properties of the function $s(\xi)$ dictated by the causality principle are such that the remaining factor $\left(\xi-c_{0} V^{2} s(\xi)\right)^{-1}$ has no singularities in the lower part of the complex plane, and therefore by transforming the initial integration contour along the real axis into the integration contour $C_{2}$, shown in Fig. 3.2 (the circumvention direction is the opposite to that shown in the figure) one gets the numerical
contribution $1 / 2$ implied by the condition $s(\xi) \rightarrow 0$ at $|\xi| \rightarrow \infty$. The other half comes from the complex conjugate (the second) term in (4.181).

The equation for the population difference is similar to (4.29). It emerges from (4.176) (4.177) with the allowance of (4.180) and reads

$$
\begin{align*}
\Delta \rho(\varepsilon, \xi)= & \left(\frac{\widetilde{F}(\varepsilon, \xi)-1}{\widetilde{F}(\varepsilon, \xi)+1} \frac{V^{2} c_{0}(s(\varepsilon)-s(\xi))}{(\varepsilon-\xi)}+1\right) \\
& \frac{1}{\xi-c_{0} V^{2} s(\xi)} \frac{1}{\varepsilon-c_{0} V^{2} s(\varepsilon)} \tag{4.182}
\end{align*}
$$

whereas the distribution over the second band results from Eqs.(4.174, 4.176) and (4.180), that yield

$$
\begin{align*}
\rho_{m}(\varepsilon, \xi)= & \frac{1}{\varepsilon-\Delta_{m}-c_{m} V^{2} s(\varepsilon)} \frac{1}{\xi-\Delta_{m}-c_{m} V^{2} s(\xi)} \\
& \frac{(s(\varepsilon)-s(\xi)) c_{0} c_{m} V^{4}}{F(\varepsilon, \xi)(\varepsilon-\xi)\left(\xi-c_{0} V^{2} s(\xi)\right)\left(\varepsilon-c_{0} V^{2} s(\varepsilon)\right)} \tag{4.183}
\end{align*}
$$

In the long time limit, the main contribution comes from the point $\varepsilon=\xi$, where $\widetilde{F}(\varepsilon, \xi)=1$ according to (4.180), and the function $F(\varepsilon, \xi)$ is regular, while the difference $s(\varepsilon)-s(\xi)$ assumes an imaginary value. Therefore (4.183) after the inverse Fourier transformation takes the form

$$
\begin{align*}
\rho_{m}(\Delta)= & \int \frac{1}{\varepsilon-\Delta-c(\Delta) V^{2} s(\varepsilon)} \frac{1}{\varepsilon-\Delta-c(\Delta) V^{2} s^{*}(\varepsilon)} \\
& \frac{2 \pi i\left(s(\varepsilon)-s^{*}(\varepsilon)\right) c_{0} c(\Delta) V^{4}}{F(\varepsilon, \varepsilon)\left(\varepsilon-c_{0} V^{2} s^{*}(\varepsilon)\right)\left(\varepsilon-c_{0} V^{2} s(\varepsilon)\right)} d \varepsilon \tag{4.184}
\end{align*}
$$

with $s^{*}(\varepsilon)=\left.s(\xi)\right|_{\xi=\varepsilon}$. We also note, that the asymptotic population distribution over the states of the first band at $t \rightarrow \infty$ coincides with that for the second band.

We now have to determine the functions $s(\varepsilon), s(\xi), F(\varepsilon, \xi)$, and $\widetilde{F}(\varepsilon, \xi)$, entering (4.183,4.184). For the model of two continuous bands of an identical state density $g$ they are given by the equations

$$
\begin{align*}
s(\varepsilon) & =\int \frac{g d \Delta}{(\varepsilon-\Delta) / c(\Delta)-V^{2} s(\varepsilon)}  \tag{4.185}\\
F(\varepsilon, \xi) & =\int \frac{V}{\frac{\varepsilon-\Delta}{c(\Delta)}-V^{2} s(\varepsilon)} \frac{V}{\frac{\xi-\Delta}{c(\Delta)}-V^{2} s(\xi)} \frac{g d \Delta}{c(\Delta)}  \tag{4.186}\\
\widetilde{F}(\varepsilon, \xi) & =\int \frac{V}{\frac{\varepsilon-\Delta}{c(\Delta)}-V^{2} s(\varepsilon)} \frac{V}{\frac{\xi-\Delta}{c(\Delta)}-V^{2} s(\xi)} g d \Delta .
\end{align*}
$$

In order to proceed with the calculations we have to choose a model for the dependence $c(\Delta)$ which would allow us to evaluate the integrals $(4.185,4.186)$
and solve the equation (4.185) for $s(\varepsilon)$ choosing for $s(\xi)$, the complex conjugate solution. However, even for the simplest model of coupling $c(\Delta)=$ $\Gamma / g \pi\left(\Delta^{2}+\Gamma^{2}\right)$ which has a Lorentzian profile of width $\Gamma$, the calculations turn out to be rather cumbersome. We therefore will show here just the main steps of the general analysis and concentrate mainly on the case of strong coupling $V^{2} g \gg \Gamma$, which can be performed analytically. We note that the opposite extreme $V^{2} g \ll \Gamma$ corresponds to the case of two homogeneous bands that has already been considered in the section 4.2 .

Replacements $\Delta \rightarrow \Gamma x, \varepsilon \rightarrow \Gamma \epsilon, s(\epsilon) \rightarrow 2 \pi g \Gamma^{2} r(\epsilon) / V^{2}, \kappa \rightarrow V^{2} / \pi^{2} g \Gamma^{3}$ in (4.185) yield

$$
\begin{equation*}
r(\epsilon)=\frac{\kappa}{2} \int \frac{d x}{(\epsilon-x)\left(x^{2}+1\right)-2 r(\epsilon)} . \tag{4.187}
\end{equation*}
$$

The denominator of the integrand is a third-order polynomial with respect to $x$, which has roots at the points

$$
\begin{align*}
& a=\frac{\epsilon}{3}+\frac{3-\epsilon^{2}-\Theta^{2 / 3}}{3 \Theta^{1 / 3}} \\
& b=\frac{\epsilon}{3}+\frac{e^{i \pi / 3}\left(\epsilon^{2}-3\right)+e^{-i \pi / 3} \Theta^{2 / 3}}{3 \Theta^{1 / 3}}, \\
& c=\frac{\epsilon}{3}+\frac{e^{-i \pi / 3}\left(\epsilon^{2}-3\right)+e^{i \pi / 3} \Theta^{2 / 3}}{3 \Theta^{1 / 3}}, \tag{4.188}
\end{align*}
$$

where

$$
\begin{equation*}
\Theta(\epsilon, r)=27 r-9 \epsilon-\epsilon^{3}+3 \sqrt{3} \sqrt{27 r^{2}+\left(1+\epsilon^{2}\right)^{2}-2 r \epsilon\left(9+\epsilon^{2}\right)} \tag{4.189}
\end{equation*}
$$

For $\operatorname{Im} \epsilon>0$ the points $a$ and $b$ are in the upper part of the complex plane, whereas the point $c$ is in the lower part. Therefore the integral (4.187) can be evaluated by taking the residue at the point $x=c$ which yields the algebraic equation

$$
\begin{equation*}
r(\epsilon)=\frac{-\pi i \kappa}{(c-a)(c-b)} . \tag{4.190}
\end{equation*}
$$

Let us make the replacement $\Theta \rightarrow\left(e^{2 i \arccos (j)-2 \pi i / 3}\left(\epsilon^{2}-3\right)\right)^{3 / 2}$, which yields

$$
\begin{align*}
& a=\frac{\epsilon}{3}-\frac{1}{3}\left(j+\sqrt{3} \sqrt{1-j^{2}}\right) \sqrt{\epsilon^{2}-3}, \\
& b=\frac{\epsilon}{3}-\frac{1}{3}\left(j-\sqrt{3} \sqrt{1-j^{2}}\right) \sqrt{\epsilon^{2}-3}, \\
& c=\frac{\epsilon}{3}+\frac{2}{3} j \sqrt{\epsilon^{2}-3}, \tag{4.191}
\end{align*}
$$

for the roots given by (4.188) and thereby allows one to write the expressions (4.189) and (4.190) in a much shorter form

$$
\begin{align*}
& r(\epsilon)=\frac{2\left(\epsilon^{2}-3\right)^{3 / 2}}{27}\left[4 j^{3}-3 j+\frac{\epsilon\left(\epsilon^{2}+9\right)}{\left(\epsilon^{2}-3\right)^{3 / 2}}\right]  \tag{4.192}\\
& r(\epsilon)=\frac{3 \pi i \kappa}{\left(\epsilon^{2}-3\right)\left(4 j^{2}-1\right)} . \tag{4.193}
\end{align*}
$$

One sees that this system of equations gives an explicit expression for $r(\epsilon)$ provided one has selected the correct branch $j(\epsilon, \kappa)$ from the five possible solutions of the 5 -th order polynomial equation

$$
\begin{equation*}
\frac{81 \pi i \kappa}{2\left(\epsilon^{2}-3\right)^{5 / 2}}=\left[4 j^{3}-3 j+\frac{\epsilon\left(\epsilon^{2}+9\right)}{\left(\epsilon^{2}-3\right)^{3 / 2}}\right]\left(4 j^{2}-1\right) \tag{4.194}
\end{equation*}
$$

and substitutes this solution into (4.192). The correct branch has the asymptotic form $j(\epsilon, \kappa) \rightarrow 1$ for $\epsilon \rightarrow \infty$, as it follows from the direct evaluation of the integral (4.187) in the limit

$$
\begin{equation*}
r(\epsilon) \simeq \frac{\pi i \kappa}{\epsilon^{2}} \tag{4.195}
\end{equation*}
$$

when the main contribution comes from the pole $x=\epsilon$.
The strong coupling limit $V^{2} g \gg \Gamma$ corresponds to the extreme $\kappa \rightarrow \infty$, where apart from (4.195), which is valid in the domain $\Delta / \Gamma \sim \epsilon \sim \kappa^{1 / 3} \gg 1$ containing the main part of the population distribution, the approximate expressions

$$
\begin{align*}
& j \simeq 1,  \tag{4.196}\\
& c \simeq \frac{\epsilon}{3}+\frac{2}{3} \sqrt{\epsilon^{2}-3}
\end{align*}
$$

can also be found for the auxiliary variable $j$ and for the root $c$. The main contributions to the integrals of (4.185)-(4.186) and to the inverse Fourier transformation come from the point $\xi \simeq \varepsilon \simeq \Delta, \quad|\varepsilon-\xi| \ll|\varepsilon|$. With allowance of the relations $c_{0}=1 / g \pi \Gamma, c(\Delta) \simeq \Gamma / g \pi \Delta^{2}$ this yields

$$
\begin{align*}
s(\varepsilon) & \simeq \frac{2 i \Gamma}{\varepsilon^{2}} \\
\widetilde{F}(\varepsilon, \xi) & \simeq \frac{2 i \Gamma^{2} V^{2}\left(\varepsilon^{-4}+\xi^{-4}\right)}{g \pi(\xi-\varepsilon)+2 i \Gamma^{2} V^{2}\left(\varepsilon^{-4}+\xi^{-4}\right)} \tag{4.197}
\end{align*}
$$

Substitution of (4.197) (4.182) gives the population difference

$$
\begin{align*}
\Delta \rho(\varepsilon, \xi)= & \left(\frac{-V^{2} 4 i \varepsilon^{2}}{8 i \Gamma^{2} V^{2}+g \pi \varepsilon^{4}(\xi-\varepsilon)}+1\right) \\
& \frac{\varepsilon^{2} g \pi}{\varepsilon^{3} g \pi-2 V^{2} i} \frac{\xi^{2} g \pi}{\xi^{3} g \pi+2 V^{2} i} \tag{4.198}
\end{align*}
$$



Fig. 4.26. Two identical bands with mean squared interaction depending on the position of the coupled levels: $W_{n, m}\left(\Delta_{n}, \Delta_{m}\right)=\Gamma^{4}\left\langle V^{2}\right\rangle g /\left[\left(\Gamma^{2}+\Delta_{m}^{2}\right)\left(\Gamma^{2}+\Delta_{n}^{2}\right)\right]$ in the regime $\Gamma \ll\left\langle V^{2}\right\rangle g$. (a) The difference of the total populations of the first and the second bands depends on the scaled time $T=2 \pi t \Gamma^{2} g /\left\langle V^{2}\right\rangle$. (b) The asymptotic distribution of the populations over the energies of the band states is identical for the first and the second bands and depends only on the scaled detuning $\Delta / a$, where $a=\left(2\left\langle V^{2}\right\rangle / g \pi\right)^{1 / 3}$.

We now perform the inverse Fourier transformation of (4.198) and obtain the time dependent population $\Delta \rho(t)$ shown in Fig. 4.26(a).

In order to obtain this dependence we consider first the contribution of the pole at $\xi=i\left(2 V^{2} / g \pi\right)^{1 / 3}$ in the upper part of the complex plane $\xi$ where the integration contour can be displaced for positive $t$. For this pole the only singularity in the lower part of the complex plane of $\varepsilon$ is located at $\varepsilon=-i\left(2 V^{2} / g \pi\right)^{1 / 3}$, and even at $t=0$ it yields a small $\sim \kappa^{-2 / 3}$ contribution to the expression (4.198), which becomes yet smaller, decreasing exponentially with elapsing time. It can therefore be completely ignored for the given order of approximation. The main contribution to the inverse Fourier transformation for the population difference comes from the point $\varepsilon \simeq \xi$ and reads

$$
\begin{equation*}
\Delta \rho(t)=\int_{C} e^{-t 8 \Gamma^{2} V^{2} / g \pi \varepsilon^{4}} \frac{2 g V^{2} \varepsilon^{2}}{\varepsilon^{6}(g \pi)^{2}+4 V^{4}} d \varepsilon . \tag{4.199}
\end{equation*}
$$

The replacements $\varepsilon^{3} \rightarrow 2 x V^{2} / g \pi$ and $t \rightarrow T\left(2 V^{2} / g \pi\right)^{1 / 3} / 4 \Gamma^{2}$ yield

$$
\begin{align*}
\Delta \rho(t) & =\frac{1}{\pi} \int_{C_{x}} e^{-T x^{-4 / 3}} \frac{d x}{3\left(x^{2}+1\right)} \\
& =\frac{3}{\pi} \int_{-\infty}^{\infty} e^{-T x^{-4 / 3}} \frac{d x}{3\left(x^{2}+1\right)}, \tag{4.200}
\end{align*}
$$

where the contour $C_{x}$ circumvents the point $x=0$ three times, in order to accumulate the same phase difference of $6 \pi$ as the function $\varepsilon^{3}$ after a single circumvention of the point $\varepsilon=0$ along the contour $C$. This is the origin of the topological factor 3 in the last equality of (4.200). At $t=0$ this equation gives $\Delta \rho(0)=1$, and for $t \rightarrow \infty$ it yields the asymptotic power-law dependence

$$
\begin{equation*}
\Delta \rho(t)=\frac{6 V^{2}}{g \pi^{2}}\left(\frac{g \pi}{8 \Gamma^{2} V^{2} t}\right)^{3 / 4} \int_{-\infty}^{\infty} e^{-1 / u^{4}} \frac{d u}{u^{4}} \sim t^{-3 / 4} \tag{4.201}
\end{equation*}
$$

Let us consider now the population distribution over the energy levels of the bands suggested by (4.184). With the allowance of (4.180)-(4.197), the relations $c_{0}=1 / g \pi \Gamma, c(\Delta) \simeq \Gamma / g \pi \Delta^{2}$, and the topological factor 3 mentioned above it reads

$$
\begin{equation*}
\rho_{m}(\Delta)=\frac{6 V^{2}}{(g \pi)^{2}} \frac{\Delta^{2}}{\Delta^{6}+\left(\frac{2 V^{2}}{g \pi}\right)^{2}} \tag{4.202}
\end{equation*}
$$

which corresponds to a two-hump distribution of width $\left(2 V^{2} / g \pi\right)^{1 / 3}$, shown in Fig.4.26(b). We note that $\int \rho_{m}(\Delta) g d \Delta=1 / 2$, since the population is equally distributed amongst the states of the first and the second band of the same energy. We also note that for the regime $V^{2} g \gg \Gamma$ under consideration the width $\Gamma$ of the coupling line $W(\Delta)$ does not effect the population distribution profile (4.202), which is located on the far wings of this line.

## 5 Soluble Time-Dependent Systems

We now consider quantum systems with time-dependent Hamiltonians. The formal expression for the evolution operator for such systems is given by the time-ordered exponent of the integral of the corresponding Hamiltonian

$$
\begin{align*}
\widehat{U}(t)= & \widehat{T} \exp \left\{-i \int^{t} \widehat{H}(\tau) d \tau\right\}=e^{-i \widehat{H}(t) d \tau} e^{-i \widehat{H}(t-d \tau) d \tau} \ldots e^{-i \widehat{H}(\tau+d \tau) d \tau} \\
& e^{-i \widehat{H}(\tau) d \tau} \ldots e^{-i \widehat{H}(2 d \tau) d \tau} e^{-i \widehat{H}(d \tau) d \tau} e^{-i \widehat{H}(0) d \tau} \tag{5.1}
\end{align*}
$$

which simply involves the sequential application of infinitesimal evolution operators $e^{-i \widehat{H}(\tau) d \tau}$ with instantaneous Hamiltonians. Here $\widehat{T}$ denotes the time ordering operator. Usually this expression is not very practical when one looks for an exact analytic solution. Moreover, it turns out that for a multilevel quantum system, the possibility to achieve an analytic description is strongly related to the algebraic structure of the time-dependent Hamiltonian suggested by Lie group theory. We therefore start our consideration by introducing some general ideas, remarks, and results of the group theory required for our analysis.

We continue with the analysis of a number of exactly soluble problems widely known as the Landau-Zener problem and its generalization to the case of decaying and dephasing levels. We also consider a two-level system for two special cases of the time-dependent perturbation that can be performed analytically. The results are obtained in terms of the hypergeometric functions which is almost always the case for the exactly soluble problems. In certain conditions, the hypergeometric functions can be reduced to other special functions, which is the case for exponentially rising coupling. Then we turn to the adiabatic approximation and show that the probability of nonadiabatic transitions is given by the Dykhne formula, also valid for multilevel time-dependent systems.

The next classical example of multilevel time-dependent systems is the Demkov-Osherov system consisting of a level moving across a band. This problem can be solved exactly by the Laplace contour integral method. In this context we address a question about non-adiabatic processes at the edges of the band. The Laplace method allows one to find a solution in the form of a contour integral for another problem - a level placed at a certain energy difference underneath a continuous band with the coupling exponentially in-
creasing between the level and the band. It turns out that in spite of the fact that with the course of time, the coupling rapidly exceeds the detuning between the level and the edge of the band, the level keeps the main part of its population, being only repelled toward negative energies.

### 5.1 Algebraic Structure of Time Dependent Systems

The dynamics of a quantum system with time-dependent parameters is often more involved compared to the stationary analog. This is always the case when the system Hamiltonians corresponding to different moments in time do not commute. The algebra of their commutators of various orders, namely the rules by which the commutators can be expressed in terms of a linear combinations of operators $\widehat{\mathcal{G}}_{n}$ (where $n=0 \ldots N_{o}$, and $\widehat{\mathcal{G}}_{0}=\widehat{I}$ ) comprising a set of $N_{o}$ so-called generators of a Lie group, play an important role and shift the consideration of the problem to the domain of Lie group theory. In the following sections, we illustrate the group theory background of several exactly soluble problems, and here we demonstrate it for the simplest example of a system with a Hamiltonian that takes only two values $\widehat{H}=\widehat{H}_{A}$ or $\widehat{H}=\widehat{H}_{B}$.

Consider the system after it has been evolving with the Hamiltonian $\widehat{H}_{A}$ during a time interval $t_{A}$ and with the Hamiltonian $\widehat{H}_{B}$ during a time interval $t_{B}$. The evolution operator reads

$$
\begin{equation*}
\widehat{U}\left(t_{A}, t_{B}\right)=e^{-i t_{B} \widehat{H}_{B}} e^{-i t_{A} \widehat{H}_{A}} . \tag{5.2}
\end{equation*}
$$

The essence of the problem is in the fact that the regular multiplication rule $e^{B} e^{A}=e^{A+B}$ is not valid for non-commuting operators $\widehat{A}$ and $\widehat{B}$, and must be replaced by the relation known as the Baker-Campbell-Hausdorff formula

$$
\begin{gather*}
e^{\widehat{B}} e^{\widehat{A}}=e^{\widehat{C}}, \\
\widehat{C}=\widehat{A}+\widehat{B}+\frac{1}{2}[\widehat{A}, \widehat{B}]+\frac{1}{12}[\widehat{A},[\widehat{A}, \widehat{B}]] \\
 \tag{5.3}\\
-\frac{1}{12}[\widehat{B},[\widehat{A}, \widehat{B}]]+\ldots+c_{k, s}[A,[\ldots,[A, B] \ldots]]+\ldots,
\end{gather*}
$$

where an explicit albeit non-trivial computation rule exists (see p. 220) for the numerical coefficients $c_{k, s}$ in front of a generic commutator of $k$-th order. The index $s=1, \ldots, 2^{k}$ enumerates different sequences of commutation and has unities and zeros in its binary representation in the order suggested by the sequence of operators $\widehat{A}$ and $\widehat{B}$ respectively in the corresponding multiple commutator. For example $[A,[A,[B,[A, B]]]]$ corresponds to $s=11010$ in the binary representation, that is $s=0+1 \times 2^{1}+0 \times 2^{2}+1 \times 2^{3}+1 \times 2^{4}=$ 26. Now one sees that, indeed, the commutation relations generated by two Hamiltonians $\widehat{H}_{A}$ and $\widehat{H}_{B}$ play a crucial role in the dynamics of the quantum system since

$$
\begin{align*}
U\left(t_{A}, t_{B}\right)= & \exp \left\{-i t_{B} \widehat{H}_{B}-i t_{A} \widehat{H}_{A}-\frac{t_{B} t_{A}}{2}\left[\widehat{H}_{B}, \widehat{H}_{A}\right]\right. \\
& \left.-\frac{i t_{B}^{2} t_{A}}{12}\left[\widehat{H}_{B},\left[\widehat{H}_{B}, \widehat{H}_{A}\right]\right]+\frac{i t_{B} t_{A}^{2}}{12}\left[\widehat{H}_{A},\left[\widehat{H}_{B}, \widehat{H}_{A}\right]\right]+\ldots\right\} \tag{5.4}
\end{align*}
$$

Equation (5.4) gives us a clue for the classification of different particular cases - the larger the dimension of the operator space produced by the commutators, that is the larger the set of operators comprising the algebra, the more involved the quantum dynamics are. For a commuting pair of Hamiltonians the set has only two elements since the evolution operator is given by the exponent of the linear combination $-i t_{B} \widehat{H}_{B}-i t_{A} \widehat{H}_{A}$ of just two operators, and hence the operator space is two-dimensional. One can say that the evolution operator (5.2) belongs to a two-parametric Lie group. Apart from a phase factor $e^{-\frac{t_{B} t_{A}}{2} C}$ it is also the case when the commutator is a number $C=\left[\widehat{H}_{B}, \widehat{H}_{A}\right]$. In a more involved case, when all of the commutators of the order higher than a certain $k$ vanish, the evolutionary operator turns to be the exponent of a linear combination of a larger, albeit finite number of terms, each of which is a commutator of some order less than $k$. The dimensionality $N_{o}$ of the operator space can reach at most $2^{k}$, since this value gives the number of different commutators, unless some of the commutators can be represented as linear combinations of the others, and hence the number of linearly independent elements of the set comprising the algebra is smaller.

Another case, the most interesting from the viewpoint of a comprehensive analytical consideration, corresponds to the situation when no commutators of higher order vanish, but can all be represented as a linear combination of commutators of orders less than a certain order $k$. In this case, all of the coefficients in (5.3) are different from zero, but nevertheless this infinite sum can be cast in a finite linear combination of commutators with the coefficients $f_{n}\left(t_{A}, t_{B}\right)$ given by the infinite series. In this very case, Lie group theory is the most important concept, which allows one to employ generators $\mathcal{G}_{n}$ of standard groups as a basis set in the space of commutators and thereby express the evolution operator in the form

$$
\begin{equation*}
\widehat{U}(t)=e^{\sum_{n=0}^{N_{o}} f_{n}\left(t_{A}, t_{B}\right) \widehat{\mathcal{G}}_{n}} \tag{5.5}
\end{equation*}
$$

which sometimes provides one with an exact analytical solution of the problem. In other words, this is the case where the dimensionality of the operator space generated by the commutators of the Hamiltonians at different moments of time is much less than the dimensionality of the Hilbert space of the considered system: $N_{o} \ll N$. The evolution operator belongs to an $N_{o^{-}}$ parametric Lie group, although the values of these parameters $f_{n}\left(t_{A}, t_{B}\right)$ are possibly rather complex functions of $t_{A}$ and $t_{B}$ even for the simplest particular problem (5.2) under consideration. Finally, in the case where all of the commutators of the lower orders turn out to be linearly independent until their total number $N_{o}$ reaches the maximum possible number $N^{2}$ of linearly independent Hermitian matrices in a $N$-dimensional Hilbert space of the system, we encounter the so-called non-holonom problem. No compact general
solution of the problem in terms of analytic or special functions exists for large $N$, and the analytic description requires statistical hypotheses, as is the case for complex stationary quantum systems.

In order to consider time-dependent Hamiltonians of a general type one could consider time as one of the parameters of the Lie group and treat the operator of the time derivative as one of the group generators. From the formal point of view this consideration relies on the method of characteristics for linear partial deferential equations as well as on the operator of a finite shift. We will not however dwell on the mathematical background of the approach but instead illustrate this formal construction by the following physical example. Consider a point quantum system with a time-dependent Hamiltonian $\widehat{H}(t)$ for which the Shrödinger equation reads

$$
\begin{equation*}
i|\dot{\Psi}(t)\rangle=\widehat{H}(t)|\Psi(t)\rangle \tag{5.6}
\end{equation*}
$$

Imagine that the time dependence of the Hamiltonian originates from the fact that the system moves with a constant velocity $v$ along the axis $x$, whereas the Hamiltonian $\widehat{H}(t)=\left.\widehat{H}(x / v)\right|_{x=v t}$ depends on the position $x(t)=v t$ of the system on this axis. Equation (5.6) describes the evolution of the state vector $|\Psi(t)\rangle$ in the reference frame attached to the moving system, whereas in the reference frame at rest, the Schrödinger equation adopts the form

$$
\begin{equation*}
i|\dot{\Psi}(t, x)\rangle=\left[\widehat{H}\left(\frac{x}{v}\right)-i v \frac{\partial}{\partial x}\right]|\Psi(t, x)\rangle, \tag{5.7}
\end{equation*}
$$

where the last term in the square brackets allows for the Galilean transformation to a reference frame moving with a constant velocity. Once the solution $|\Psi(t, x)\rangle$ of (5.7) is found, one should set $x=v t$. To shorten the notation, we take $x=v \tau$ and rewrite (5.7) in the form

$$
\begin{equation*}
i|\dot{\Psi}(t, \tau)\rangle=\left[\widehat{H}(\tau)-i \frac{\partial}{\partial \tau}\right]|\Psi(t, \tau)\rangle \tag{5.8}
\end{equation*}
$$

where the new effective Hamiltonian $\left[\hat{H}(\tau)-i \frac{\partial}{\partial \tau}\right]$ depends on the new "time" variable $\tau$, but does not depend on "real" time $t$. Therefore the operator of evolution corresponding to (5.8) reads

$$
\begin{equation*}
\widehat{U}(t)=e^{-i t\left[\widehat{H}(\tau)-i \frac{\partial}{\partial \tau}\right]} \tag{5.9}
\end{equation*}
$$

which implies

$$
\begin{equation*}
|\Psi(t, \tau)\rangle=e^{-i t\left[\widehat{H}(\tau)-i \frac{\partial}{\partial \tau}\right]}|\Psi(0, \tau)\rangle \tag{5.10}
\end{equation*}
$$

given an initial condition $|\Psi(0, \tau)\rangle$ at $t=0$. One can choose an arbitrary initial condition for the state vector $|\Psi(t, \tau)\rangle$ in two-dimensional time space, which is restricted only by the correspondence requirement $|\Psi(0,0)\rangle=|\Psi(0)\rangle$ suggested by the initial condition for the original state vector $|\Psi(t)\rangle$ entering (5.6).

Let us consider first the evolution operator (5.9) in the particular case $\widehat{H}(\tau)=0$, that is

$$
\begin{equation*}
\widehat{U}(t)=e^{-t \frac{\partial}{\partial \tau}} \tag{5.11}
\end{equation*}
$$

The application of this operator to an arbitrary function or a state vector, say $|\Psi(0, \tau)\rangle$, with allowance for a Taylor expansion formula yields

$$
\begin{equation*}
e^{-t \frac{\partial}{\partial \tau}}|\Psi(0, \tau)\rangle=\sum_{n=0}^{\infty} \frac{(-t)^{n}}{n!}\left(\frac{\partial}{\partial \tau}\right)^{n}|\Psi(0, \tau)\rangle=|\Psi(0, \tau-t)\rangle . \tag{5.12}
\end{equation*}
$$

One sees that application of the operator (5.11) results in a shift of the argument $\tau$ of an arbitrary function or state vector by a finite value $-t$. Therefore one can refer to $e^{t \frac{\partial}{\partial \tau}}$ as the operator of finite shift by the value $t$. Note the evident consequence of $(5.10)$ for $\widehat{H}(\tau)=0$ and (5.12):

$$
\begin{equation*}
|\Psi(t, \tau)\rangle=e^{-t \frac{\partial}{\partial \tau}}|\Psi(0, \tau)\rangle=|\Psi(0, \tau-t)\rangle, \tag{5.13}
\end{equation*}
$$

therefore

$$
\begin{equation*}
|\Psi(t)\rangle=\left.|\Psi(t, \tau)\rangle\right|_{\tau=t}=\left.|\Psi(0, \tau-t)\rangle\right|_{\tau=t}=|\Psi(0,0)\rangle=|\Psi(0)\rangle \tag{5.14}
\end{equation*}
$$

as it should be for $\widehat{H}(\tau)=0$, which simply indicates displacement of a nonevolving system in space with constant velocity.

Let us now consider the evolution operator (5.9) in the general case, rewrite it in the equivalent form

$$
\begin{equation*}
\widehat{U}(t)=e^{-t \frac{\partial}{\partial \tau}} e^{t \frac{\partial}{\partial \tau}} e^{-i t\left[\widehat{H}(\tau)-i \frac{\partial}{\partial \tau}\right]} \tag{5.15}
\end{equation*}
$$

and apply the Baker-Campbell-Hausdorff formula (5.3) to the last two factors. This yields

$$
\begin{equation*}
\widehat{U}(t, \tau)=e^{-t \frac{\partial}{\partial \tau}} e^{-i t \widehat{H}(\tau)-t \frac{\partial}{\partial \tau}+t \frac{\partial}{\partial \tau}+\frac{1}{2}\left[-i t \widehat{H}(\tau)-t \frac{\partial}{\partial \tau}, t \frac{\partial}{\partial \tau}\right]+\ldots} . \tag{5.16}
\end{equation*}
$$

One sees that the operator of finite shift $e^{-t \frac{\partial}{\partial \tau}}$ acts on the exponent of the argument

$$
\begin{gather*}
-i t \widehat{H}(\tau)+\frac{1}{2}\left[-i t \widehat{H}(\tau), t \frac{\partial}{\partial \tau}\right]+\frac{1}{12}\left[t \frac{\partial}{\partial \tau},\left[-i t \widehat{H}(\tau), t \frac{\partial}{\partial \tau}\right]\right]  \tag{5.17}\\
-\frac{1}{12}\left[-i t \widehat{H}(\tau)-t \frac{\partial}{\partial \tau},\left[-i t \widehat{H}(\tau), t \frac{\partial}{\partial \tau}\right]\right]+\ldots
\end{gather*}
$$

which, apart from the time dependent Hamiltonian $\widehat{H}(\tau)$ comprises commutators of different orders of the Hamiltonian and the time derivative $\frac{\partial}{\partial \tau}$. From the relation

$$
\begin{align*}
& {\left[\frac{\partial}{\partial \tau}, \widehat{H}(\tau)\right]|\Psi(t, \tau)\rangle=\frac{\partial}{\partial \tau} \widehat{H}(\tau)|\Psi(t, \tau)\rangle-\widehat{H}(\tau) \frac{\partial}{\partial \tau}|\Psi(t, \tau)\rangle} \\
& =\left[\frac{\partial}{\partial \tau} \widehat{H}(\tau)\right]|\Psi(t, \tau)\rangle+\widehat{H}(\tau) \frac{\partial}{\partial \tau}|\Psi(t, \tau)\rangle-\widehat{H}(\tau) \frac{\partial}{\partial \tau}|\Psi(t, \tau)\rangle \\
& =\widehat{\dot{H}}(\tau)|\Psi(t, \tau)\rangle \tag{5.18}
\end{align*}
$$

it becomes evident that commutation of a time dependent operator with the time derivative yields the time derivative of this operator. Equation (5.16) now reads

$$
\begin{align*}
\widehat{U}(t, \tau)= & e^{-t \frac{\partial}{\partial \tau}} \exp \left\{-i t \widehat{H}(\tau)-\frac{i t^{2}}{2} \widehat{\dot{H}}(\tau)+\ldots\right.  \tag{5.19}\\
& +\frac{i t^{k} a_{k, n, \ldots, m, l}}{2}\left[\widehat{H}^{(n)}(\tau),\left[\ldots,\left[\widehat{H}^{(m)}(\tau), \widehat{H}^{(l)}(\tau)\right] \ldots\right]+\ldots\right\} .
\end{align*}
$$

where dots denote commutators of different orders of the time derivatives $\dot{H}(\tau), \widehat{H}(\tau), \ldots, \widehat{H}^{(n)}(\tau) \ldots$ of the Hamiltonian that enter the series with the coefficients $a_{k, n, \ldots, m, l}$ related in a certain way to the coefficients $c_{k, s}$ in the Baker-Campbell-Hausdorff series (5.3).

Note that the approach based on the expansion (5.19) gives us a tool which helps us to determine whether or not a quantum system possesses a hidden symmetry, which would allow one to reduce the original problem to a problem of smaller dimension. That is, it allows one to identify the situation when the space spanned by the commutators of $\widehat{H}^{(n)}(\tau)$ has a smaller dimension compared to the dimension of the entire Hilbert space. Unfortunately this method does not suggest any practical way to calculate the operator of evolution, since the coefficients $c_{k, s}$ and hence $a_{k, n, \ldots, m, l}$ are given by a recurrent computation rule and there is no explicit expression available for these coefficients at the moment.

We also note that the application of the operator of finite shift in (5.16) results in

$$
\begin{equation*}
\widehat{U}(t, \tau)=e^{-i t \widehat{H}(\tau-t)-\frac{i t^{2}}{2} \widehat{\dot{H}}(\tau-t)+\ldots} \tag{5.20}
\end{equation*}
$$

and once we set $\tau=x / v=t$ (that is we remain on the characteristics, or in other words on the trajectory, see (5.6)-(5.7)) it yields the equation

$$
\begin{align*}
\widehat{U}(t)= & \exp \left\{-i t \widehat{H}(0)-\frac{i t^{2}}{2} \widehat{\dot{H}}(0)-\frac{i t^{3}}{3!} \widehat{\ddot{H}}(0)+\ldots\right. \\
& +a_{k, n, \ldots, m, l} \frac{i t^{k}}{2}\left[\widehat{H}^{(n)}(0),\left[\ldots,\left[\widehat{H}^{(m)}(0), \widehat{H}^{(l)}(0)\right] \ldots\right]+\ldots\right\} \\
= & \widehat{T} \exp \left\{-i \int^{t} \widehat{H}(\tau) d \tau\right\} \tag{5.21}
\end{align*}
$$

which plays the role of the Taylor expansion for the integral in the timeordered exponent $\widehat{T} \exp \left\{\int^{t} \widehat{H}(\tau) d \tau\right\}$ encountered earlier in (5.1). Indeed, in the first line of (5.21) one recognizes the Taylor expansion of the integral $\int^{t} \widehat{H}(\tau) d \tau$ such that (5.21) has the structure

$$
\begin{align*}
\widehat{U}(t) & =\widehat{T} \exp \left\{-i \int^{t} \widehat{H}(\tau) d \tau\right\} \\
& =\exp \left[-i \int^{t} \widehat{H}(\tau) d \tau-i \widehat{R}\left(t ; \widehat{H}^{(n)}(0)\right)\right] \tag{5.22}
\end{align*}
$$

where the remainder $\widehat{R}\left(t, \widehat{H}^{(n)}(0)\right)$ is a Hermitian operator, which has a complex explicit form, contains the commutators of the Hamiltonian and its derivatives of different orders, and thereby allows for the non-commutativity of the Hamiltonian corresponding to different moments in time.

Consider an important particular case of time dependence. A Hamiltonian periodic in time of period $2 \pi / \omega$ can be represented as a Fourier series $\widehat{H}(\tau)=$ $\sum_{n} \widehat{H}_{n} e^{-i \omega n \tau}$. One can also cast the state vector into a Fourier series

$$
\begin{equation*}
|\Psi(t, \tau)\rangle=\sum_{n} e^{-i \omega n \tau}\left|\Psi_{n}(t)\right\rangle \tag{5.23}
\end{equation*}
$$

and rewrite (5.8) in the form

$$
\begin{equation*}
i\left|\dot{\Psi}_{m}(t)\right\rangle=\sum_{k} \widehat{H}_{m-k}\left|\Psi_{k}(t)\right\rangle-\omega m\left|\Psi_{m}(t)\right\rangle, \tag{5.24}
\end{equation*}
$$

which is a system of linear first-order differential equations with time independent coefficients. One can represent the state vector (5.23) as a double sum

$$
\begin{equation*}
|\Psi(t, \tau)\rangle=\sum_{n, j} e^{-i \omega n \tau} \psi_{n, j}(t)|j\rangle \tag{5.25}
\end{equation*}
$$

where the states $|j\rangle$ form a complete basis set for the system, which yields

$$
\begin{equation*}
i \dot{\psi}_{m, j}(t)=\sum_{k, l} H_{m-k, j, l} \psi_{k, l}(t)-\omega m \psi_{m, l}(t) \tag{5.26}
\end{equation*}
$$

The Fourier transformation of this equation over the time $t$ yields

$$
\begin{equation*}
E \psi_{m, l}=\sum_{k, l} H_{m-k, j, l} \psi_{k, l}-\omega m \psi_{m, l} \tag{5.27}
\end{equation*}
$$

The eigenvalues $E_{\nu}$ of this equation are called quasienergies by analogy with the quasimomenta for the systems with Hamiltonians periodic in space, while such a form of the time-dependent Schrödinger equation (5.6) is called the quasienergy representation. This method is also known as the Floquet approach.

The initial condition for (5.26) reads

$$
\begin{equation*}
|\Psi(t=0)\rangle=\sum_{n, j} \psi_{n, j}(t=0)|j\rangle \tag{5.28}
\end{equation*}
$$

Later on we will see that for the case where the time dependence results from the interaction of the system with a quantized external electromagnetic field, the coefficients $\psi_{n, j}(0)$ have the clear physical interpretation of the probability amplitudes to find the system in the quantum state $|j\rangle$ and the field in the state with $n$ photons at $t=0$. However, for a classical electromagnetic
field, or in the general case of a periodic Hamiltonian, the initial condition (5.28) alone does not permit one to determine unambiguously all of the coefficients $\psi_{n, j}(t=0)$, as was the case for (5.10). One of the most common ways to choose the initial condition for the amplitudes $\psi_{n, j}(t=0)$ is to take them as identical for all the quasienergy indices $\psi_{n, j}=\psi_{n^{\prime}, j}$, which actually corresponds to the classical limit of the coherent quantum field state.

### 5.2 Time-Dependent Two-Level Systems

The simplest example of systems with time-dependent Hamiltonians are twolevel systems. The corresponding Schrödinger equation has the form of a set of two linear, first-order differential equations

$$
\begin{align*}
i \dot{\psi}_{0} & =\Delta_{0}(t) \psi_{0}+V(t) \psi_{1} \\
i \dot{\psi}_{1} & =\Delta_{1}(t) \psi_{1}+V(t) \psi_{0} \tag{5.29}
\end{align*}
$$

for the probability amplitudes $\psi_{0}$ and $\psi_{1}$ to be in the states $|0\rangle$ and $|1\rangle$ respectively, and is similar to (3.7) for the time-independent Hamiltonian. It can also be rewritten in the equivalent form of one linear second-order differential equation with time-dependent coefficients

$$
\begin{equation*}
\ddot{\psi}_{1}+\left(\frac{\dot{V}}{V}-i \Delta_{1}+i \Delta_{0}\right) \dot{\psi}_{1}+\Delta_{1}\left(i \frac{\dot{\Delta}_{1}}{\Delta_{1}}-i \frac{\dot{V}}{V}-\Delta_{0}+\frac{V^{2}}{\Delta_{1}}\right) \psi_{1}=0 \tag{5.30}
\end{equation*}
$$

which just by interchange of indices yields the analagous equation for $\psi_{0}$. Solutions of a large number of second-order linear differential equations are known. They can be found and expressed in terms of special functions, provided a proper replacement and change of the variable equation (5.30) is reduced to one of the standard forms of the second-order equations for the special functions. In this section we consider several examples that give insight into the typical behavior of two-level systems in various conditions.

### 5.2.1 Landau-Zener Problem

In Fig. 5.1 we show an example of an exactly soluble problem, known as the Landau-Zener problem, considering a two-level system with a linear time dependence $\alpha t$ of the energy distance between the levels and a time-independent coupling $V$ of the levels. The corresponding Schrödinger equation reads

$$
\begin{align*}
& i \dot{\psi}_{0}=\frac{\alpha t}{2} \psi_{0}+V \psi_{1} \\
& i \dot{\psi}_{1}=-\frac{\alpha t}{2} \psi_{1}+V \psi_{0} \tag{5.31}
\end{align*}
$$

In order to illustrate the speculations of the previous section, it is convenient to represent the Hamiltonian


Fig. 5.1. Landau-Zener problem. (a) Two states move toward each other with linear dependence of the energy difference with time. At $t=0$ the level-crossing occurs, which means that the positions of the levels on the energy scale coincide, and for positive times the levels continue their motion with the same velocity $\alpha$. The states are coupled by the interaction $V$. (b) Position of the level energies and the Hamiltonian eigenstates of the system as a function of time. The dot-dash lines show so-called diabatic levels, that is the energy position of the states in the absence of the interaction. The dashed lines show the so-called adiabatic levels, that is the energy eigenstates of the system corresponding to the instantaneous Hamiltonian.

$$
\begin{equation*}
\widehat{H}=\frac{\alpha t}{2} \widehat{\sigma}_{z}+V \widehat{\sigma}_{x} \tag{5.32}
\end{equation*}
$$

with the help of three Pauli matrices

$$
\widehat{\sigma}_{x}=\left(\begin{array}{cc}
0 & 1  \tag{5.33}\\
1 & 0
\end{array}\right) ; \widehat{\sigma}_{y}=\left(\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right) ; \widehat{\sigma}_{z}=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)
$$

that together with the unit matrix $\widehat{I}$ generate the $S U(2)$ group and comprise a closed Lie algebra of three elements $\widehat{\mathcal{G}}_{1}=\widehat{\sigma}_{x} ; \widehat{\mathcal{G}}_{2}=\widehat{\sigma}_{y} ; \widehat{\mathcal{G}}_{3}=\widehat{\sigma}_{z}$ satisfying the commutation relations

$$
\begin{equation*}
\left[\widehat{\sigma}_{x}, \widehat{\sigma}_{y}\right]=2 i \widehat{\sigma}_{z} ;\left[\widehat{\sigma}_{z}, \widehat{\sigma}_{x}\right]=2 i \widehat{\sigma}_{y} ;\left[\widehat{\sigma}_{y}, \widehat{\sigma}_{z}\right]=2 i \widehat{\sigma}_{x} \tag{5.34}
\end{equation*}
$$

The evolution operator (5.5) for this case reads

$$
\begin{equation*}
\widehat{U}(t)=e^{f_{0}(t) \widehat{I}+f_{x}(t) \widehat{\sigma}_{x}+f_{y}(t) \widehat{\sigma}_{y}+f_{z}(t) \widehat{\sigma}_{z}} \tag{5.35}
\end{equation*}
$$

Such a representation does not result, however, in any significant simplification of the Landau-Zener problem since it implies utilization of four timedependent coefficients $f_{n}(t)$ instead of two complex-valued time-dependent
amplitudes $\psi_{0,1}(t)$. Nevertheless, we will trace this equivalence which is useful in other, more complicated cases where the number of different amplitudes (which is sometimes infinite), exceeds considerably the number of coefficients $f_{n}(t)$ needed for the description.

Let us now find the solution of the set (5.31). We express $\psi_{1}$ from the first equation, substitute it into the second equation of the set and arrive at

$$
\begin{equation*}
\left(i \frac{\partial}{\partial t}+\frac{\alpha t}{2}\right)\left(i \frac{\partial}{\partial t}-\frac{\alpha t}{2}\right) \psi_{0}=V^{2} \psi_{0} \tag{5.36}
\end{equation*}
$$

which has the form of the Schrödinger equation for an eigenfunction $\psi(x)$ of a particle of mass $1 / 2$ moving in the parabolic potential $-\alpha^{2} x^{2} / 4$ with energy $V^{2}+i \alpha / 2$, provided we have replaced time $t$ by the coordinate $x$. In the combinations $i \frac{\partial}{\partial t}+\frac{\alpha t}{2}$ and $\frac{\alpha t}{2}-i \frac{\partial}{\partial t}$ one recognizes the harmonic oscillator annihilation and creation operators, respectively, which are generators of the $s u(1,1)$ Lie algebra. Due to this fact, the solution of (5.36) can be explicitly expressed in terms of special functions. It can either be found in handbooks on special functions in the chapters dedicated to parabolic cylinder functions, or just calculated directly by substituting a Taylor series into (5.36) which after the replacement $\psi_{0}(t)=\varphi_{0}(t) e^{-i \alpha t^{2} / 4}$ takes the form

$$
\begin{equation*}
\left(\frac{\partial^{2}}{\partial t^{2}}-i \alpha t \frac{\partial}{\partial t}+V^{2}\right) \varphi_{0}(t)=0 \tag{5.37}
\end{equation*}
$$

Two linearly independent solutions correspond to even $\varphi_{0}(t)=\sum_{n=0}^{\infty} \frac{a_{n}}{n!} t^{2 n}$ and odd $\varphi_{0}(t)=\sum_{n=0}^{\infty} \frac{b_{n}}{n!} t^{2 n+1}$ functions of time. Direct substitution of these series into (5.37) yields recurrence relations

$$
\begin{gather*}
\left(n+\frac{1}{2}\right) a_{n+1}=\frac{i \alpha}{2}\left(n-\frac{V^{2}}{2 i \alpha}\right) a_{n} \\
\left(n+\frac{3}{2}\right) b_{n+1}=\frac{i \alpha}{2}\left(n-\frac{V^{2}}{2 i \alpha}+\frac{1}{2}\right) b_{n} \tag{5.38}
\end{gather*}
$$

for the coefficients $a_{n}$ and $b_{n}$. The same relations among coefficients hold for the confluent (or degenerate) hypergeometric functions

$$
\begin{equation*}
{ }_{1} F_{1}(\beta ; \gamma ; z)=1+\frac{\beta}{\gamma} \frac{z}{1!}+\frac{\beta(\beta+1)}{\gamma(\gamma+1)} \frac{z^{2}}{2!}+\ldots, \tag{5.39}
\end{equation*}
$$

when $z=i \alpha t^{2} / 2, \beta=-V^{2} / 2 i \alpha, \gamma=1 / 2$ for even and $z=i \alpha t^{2} / 2, \beta=$ $1 / 2-V^{2} / 2 i \alpha, \gamma=3 / 2$ for odd functions. Therefore, the general solution of (5.36) reads

$$
\begin{equation*}
\psi_{0}(t)=\left[c_{11} F_{1}\left(-\frac{V^{2}}{2 i \alpha} ; \frac{1}{2} ; \frac{i \alpha t^{2}}{2}\right)+c_{2} t_{1} F_{1}\left(\frac{1}{2}-\frac{V^{2}}{2 i \alpha} ; \frac{3}{2} ; \frac{i \alpha t^{2}}{2}\right)\right] e^{-i \alpha t^{2} / 4} \tag{5.40}
\end{equation*}
$$

where $c_{1}$ and $c_{2}$ are two constants. From (5.31) one sees that the replacement $\alpha \rightarrow-\alpha$ corresponds to the permutation $\psi_{0}(t) \rightleftharpoons \psi_{1}(t)$ and therefore the general solution for $\psi_{1}(t)$ reads

$$
\begin{equation*}
\psi_{1}(t)=\left[c_{31} F_{1}\left(\frac{V^{2}}{2 i \alpha} ; \frac{1}{2} ;-\frac{i \alpha t^{2}}{2}\right)+c_{4} t_{1} F_{1}\left(\frac{1}{2}+\frac{V^{2}}{2 i \alpha} ; \frac{3}{2} ;-\frac{i \alpha t^{2}}{2}\right)\right] e^{i \alpha t^{2} / 4} \tag{5.41}
\end{equation*}
$$

where the constants $c_{n}$ can be given in terms of the amplitudes at $t=0$, that is $c_{1}=\psi_{0}(0) ; c_{2}=-i V \psi_{1}(0) ; c_{3}=-i V \psi_{0}(0) ; c_{4}=\psi_{1}(0)$. In other words the evolution operator reads

$$
\widehat{U}(t)=\left(\begin{array}{cc}
{ }_{1} F_{1}\left(\frac{i V^{2}}{2 \alpha} ; \frac{1}{2} ; \frac{i \alpha t^{2}}{2}\right) e^{\frac{-i \alpha t^{2}}{4}} & \frac{V t}{i}{ }_{1} F_{1}\left(\frac{1}{2}+\frac{i V^{2}}{2 \alpha} ; \frac{3}{2} ; \frac{i \alpha t^{2}}{2}\right) e^{\frac{-i \alpha t^{2}}{4}}  \tag{5.42}\\
\frac{V t}{i}{ }_{1} F_{1}\left(\frac{1}{2}+\frac{V^{2}}{2 i \alpha} ; \frac{3}{2} ; \frac{\alpha t^{2}}{2 i}\right) e^{\frac{i \alpha t^{2}}{4}} & { }_{1} F_{1}\left(\frac{V^{2}}{2 i \alpha} ; \frac{1}{2} ; \frac{\alpha t^{2}}{2 i}\right) e^{\frac{i \alpha t^{2}}{4}}
\end{array}\right),
$$

which, after comparison with (5.35), determines implicitly the four functions $f_{n}(t)$.

One sees that the evolution operator (5.42) as well as the amplitudes (5.40)-(5.41) depend only on the dimensionless time $\sqrt{\alpha} t$ and the dimensionless adiabaticity parameter $V^{2} / \alpha$, and in Fig. 5.2 we show the population

$$
\begin{equation*}
\rho_{0}(t)=\left|\psi_{0}(t)\right|^{2}=\left|\left(\widehat{U}(t) \widehat{U}^{-1}(-\infty)\right)_{0,0}\right|^{2} \tag{5.43}
\end{equation*}
$$

of the state $|0\rangle$ as a function of these parameters, calculated for the initial condition $\psi_{0}(t=-\infty)=1$ with the help of (5.42). One can identify two limiting cases of (i) a small and of (ii) a large parameter of adiabaticity $V^{2} / \alpha$. In case (i), the coupling $V$ between the states $|0\rangle$ and $|1\rangle$ is small and the system mainly remains in its initial state following the diabatic trajectory shown in Fig. 5.1 while the population of the state $|0\rangle$ experiences just slight, dying oscillations with the phase $\sim \alpha t^{2}$ after crossing the resonance at $t=0$. In the opposite extreme (ii), the velocity $\alpha$ of the state displacement along the energy scale is small compared to $V^{2}$, such that the system follows the adiabatic trajectory and performs around $t=0$, an almost complete transition from the initial state $|0\rangle$ to the final state $|1\rangle$, although the population still manifests slight and dying oscillations corresponding to the transitions between the states.

One finds the asymptotic expression for the population

$$
\begin{equation*}
\rho_{0}(\infty)=e^{-2 \pi V^{2} / \alpha} \tag{5.44}
\end{equation*}
$$

provided the asymptotic expression for the degenerate hypergeometric function

$$
\begin{equation*}
{ }_{1} F_{1}(\beta ; \gamma ; z)=\frac{\Gamma(\gamma)}{\Gamma(\gamma-\beta)}(-z)^{-\beta} \tag{5.45}
\end{equation*}
$$

relevant to the case under consideration $(\operatorname{Re} z>0)$ and the relations

$$
\begin{align*}
|\Gamma(i y)|^{2} & =\frac{\pi}{y \sinh (\pi y)} \\
\left|\Gamma\left(\frac{1}{2}+i y\right)\right|^{2} & =\frac{\pi}{\cosh (\pi y)} \tag{5.46}
\end{align*}
$$



Fig. 5.2. (a) Population $\rho_{0}$ of the state $|0\rangle$ as a function of the dimensionless parameters $t^{2} \alpha$ and $V^{2} / \alpha$. (b) Time dependence of the population for two values of the adiabaticity parameter $V^{2} / \alpha=e^{-3}$ (upper curve) and $V^{2} / \alpha=e^{3}$ (lower curve). One sees the Rabi oscillations of the probability with the phase $\sim t^{2}$ while it tends to its asymptotic value $e^{-2 \pi V^{2} / \alpha}$ suggested by (5.44) and shown in (c)
for $\Gamma$-functions has been employed for the expression

$$
\begin{align*}
\rho_{0}(\infty)= & \lim _{t \rightarrow \infty}\left|\left(\widehat{U}(t) \widehat{U}^{-1}(-t)\right)_{0,0}\right|^{2} \\
= & \lim _{t \rightarrow \infty} \left\lvert\,{ }_{1} F_{1}\left(-\frac{V^{2}}{2 i \alpha} ; \frac{1}{2} ; \frac{i \alpha t^{2}}{2}\right)_{1} F_{1}\left(\frac{V^{2}}{2 i \alpha} ; \frac{1}{2} ;-\frac{i \alpha t^{2}}{2}\right)+\right. \\
& \left.(V t)_{1}^{2} F_{1}\left(\frac{1}{2}-\frac{V^{2}}{2 i \alpha} ; \frac{3}{2} ; \frac{i \alpha t^{2}}{2}\right)_{1} F_{1}\left(\frac{1}{2}+\frac{V^{2}}{2 i \alpha} ; \frac{3}{2} ;-\frac{i \alpha t^{2}}{2}\right)\right|^{2} \tag{5.47}
\end{align*}
$$

following from (5.43) with allowance of (5.42). The probability $\rho_{0}(\infty)=$ $e^{-2 \pi V^{2} / \alpha}$ to remain in the initial state after the level crossing is shown in Fig. 5.2(c).

From (5.44) one concludes that the probability to make a transition to the state $|1\rangle$ is $\rho_{1}(\infty)=1-e^{-2 \pi V^{2} / \alpha}$. We note that for slowly changing level positions, that is at small $\alpha$, the probability to remain in the initial state is
small, whereas the probability to make a transition is large. At the same time one can consider the problem in another frame, the so-called adiabatic basis (see Fig. 5.1), by making use of the instantaneous energy eigenstates
corresponding respectively to the eigenvalues $E_{ \pm}= \pm \sqrt{\left(\frac{\alpha t}{2}\right)^{2}+V^{2}}$ of the time-dependent Hamiltonian (5.32). Note that $|+\rangle=|0\rangle$ and $|-\rangle=|1\rangle$ at $t=-\infty$, whereas at $t=\infty$ one finds $|+\rangle=|1\rangle$ and $|-\rangle=-|0\rangle$. This means that $\rho_{+}(t=-\infty)=1$ and $\rho_{-}(t=-\infty)=0$, and the probability of making a transition from the initially populated state $|+\rangle$ to the state $|-\rangle$ is exponentially small in the limit $\alpha \rightarrow 0$, or as one says in this context adiabatically small, that is

$$
\begin{equation*}
\rho_{+}(t=\infty)=1-e^{-2 \pi V^{2} / \alpha} \tag{5.49}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{-}(t=\infty)=e^{-2 \pi V^{2} / \alpha} \tag{5.50}
\end{equation*}
$$

It is worth mentioning that the so-called Landau-Zener formula given by the expression $\rho_{d p}(\infty)=2 e^{-2 \pi V^{2} / \alpha}\left(1-e^{-2 \pi V^{2} / \alpha}\right)$ is valid for the probability of a transition following a double passage (forward and back) of the level crossing, in the absence of the quantum interference between the passages. This means that the resulting probability is just a product of the probabilities of two sequential events - to remain in the initial state after the first passage and to make a transition during the second one, while the factor 2 takes into account the possibility for the opposite order of these events. Deviation of the probability from this expression resulting from the quantum interference is known as the Ramsay effect.

### 5.2.2 Landau-Zener Transition to a Decaying State

The results of the previous section can be generalized in the case of decaying states, where the Schrödinger equation similar to (5.31) reads

$$
\begin{align*}
i \dot{\psi}_{0} & =\left(\frac{\alpha t}{2}-i \gamma_{0}\right) \psi_{0}+V \psi_{1} \\
i \dot{\psi}_{1} & =-\left(\frac{\alpha t}{2}+i \gamma_{1}\right) \psi_{1}+V \psi_{0} \tag{5.51}
\end{align*}
$$



Fig. 5.3. Landau-Zener transition between two decaying levels. The presence of the continua results in the complex shift of the time variable $\alpha t / 2 \rightarrow \alpha t / 2-i \gamma_{1 ; 2}$ where the decay rate $\gamma_{1,2}=\pi V_{0, n ; 1, m}^{2} / \delta_{0,1}$ is in accordance with (3.34).

By analogy to the level-band problem discussed in Sect.3.2 one can consider the decay rates $\gamma_{0}$ and $\gamma_{1}$ as resulting from the interaction of levels $|1\rangle$ and $|0\rangle$ with two continua of states $|n\rangle$ and $|m\rangle$ respectively, that accompany the levels in their translational motion along the energy axis, as shown in Fig. 5.3. The Schrödinger equation for this entire system does not apparently contain non-Hermitian terms proportional to $\gamma_{0}$ and $\gamma_{1}$, and indeed has a more complex form that reads

$$
\begin{align*}
i \dot{\psi}_{n} & =\left(\frac{\alpha t}{2}+n \delta_{0}\right) \psi_{n}+V_{n 0} \psi_{0} \\
i \dot{\psi}_{0} & =\frac{\alpha t}{2} \psi_{0}+V \psi_{1}+\sum_{n=-\infty}^{\infty} V_{0 n} \psi_{n} \\
i \dot{\psi}_{1} & =-\frac{\alpha t}{2} \psi_{1}+V \psi_{0}+\sum_{m=-\infty}^{\infty} V_{1 m} \psi_{m} \\
i \dot{\psi}_{m} & =\left(-\frac{\alpha t}{2}+m \delta_{1}\right) \psi_{m}+V_{m 1} \psi_{1} \tag{5.52}
\end{align*}
$$

However, since the parts of the total Hamiltonians corresponding to the interaction of the levels with their continua commute with each other as well as with the part responsible for the level coupling, the problem can actually be separated into three independent parts: the results of the levels-continua interaction yield according to (3.34) the decay rates $\gamma_{0}=\pi\left|V_{0 n}\right|^{2} / \delta_{0}$ and $\gamma_{1}=\pi\left|V_{1 m}\right|^{2} / \delta_{1}$ for the probability amplitudes, while for the interaction of the states $|1\rangle$ and $|0\rangle$ one finds (5.51). This means that the algebraic structure of the problem in the case of decaying levels remains the same as in the former case, and therefore the analytical solution can be found in terms of the same special functions. Indeed, after the replacement $\psi_{0 ; 1}(t) \rightarrow \widetilde{\psi}_{0 ; 1}(t) e^{-\left(\gamma_{0}+\gamma_{1}\right) t / 2}$
(5.51) takes the form

$$
\begin{align*}
& i \dot{\widetilde{\psi}}_{0}=\left(\frac{\alpha t-i \gamma_{0}+i \gamma_{1}}{2}\right) \widetilde{\psi}_{0}+V \widetilde{\psi}_{1} \\
& i \dot{\tilde{\psi}}_{1}=-\left(\frac{\alpha t-i \gamma_{0}+i \gamma_{1}}{2}\right) \widetilde{\psi}_{1}+V \widetilde{\psi}_{0} \tag{5.53}
\end{align*}
$$

which coincides with (5.31) for non-decaying states after the substitution $t \rightarrow t^{\prime}+i\left(\gamma_{0}-\gamma_{1}\right) / \alpha$. The evolution operator therefore reads

$$
\begin{equation*}
\widehat{\widetilde{U}}\left(t^{\prime}\right)=e^{-\left(\gamma_{0}+\gamma_{1}\right) t / 2} \widehat{U}\left(t-i \frac{\gamma_{0}-\gamma_{1}}{\alpha}\right), \tag{5.54}
\end{equation*}
$$

where the evolution operator $\widehat{U}(t)$ is given by (5.42).
The exponential pre-factor in (5.54) takes into account the exponential decrease of the total population, whereas the complex argument of the evolution operator (5.42) allows for the relative populations of the states $|1\rangle$ and $|0\rangle$ :

$$
\begin{align*}
& \rho_{0}(t) \sim\left|\left(\widehat{U}(t-i g) \widehat{U}^{-1}(-\infty-i g)\right)_{0,0}\right|^{2} \\
& \rho_{1}(t) \sim\left|\left(\widehat{U}(t-i g) \widehat{U}^{-1}(-\infty)-i g\right)_{1,0}\right|^{2} \tag{5.55}
\end{align*}
$$

where $g=\left(\gamma_{0}-\gamma_{1}\right) / \alpha$. Therefore, there is a third dimensionless parameter $\left(\gamma_{0}-\gamma_{1}\right) / \sqrt{\alpha}$ which along with dimensionless time $\sqrt{\alpha} t$ and the dimensionless adiabaticity parameter $V^{2} / \alpha$ govern the time evolution of the decaying Landau-Zener system. In Fig. 5.4 we show the dependence of the population of the state $|0\rangle$ normalized to its value $e^{-2 \gamma_{0} t}$ at $V=0$ as a function of time and the adiabaticity parameter for a positive and a negative value of $g$.

Asymptotic expressions for the degenerate hypergeometric functions

$$
\begin{equation*}
{ }_{1} F_{1}(\beta ; \gamma ; z)=\frac{\Gamma(\gamma)}{\Gamma(\beta)} z^{\beta-\gamma} e^{z} \tag{5.56}
\end{equation*}
$$

and $(5.45)$ that are valid in the case $(\operatorname{Re} z<0)$ allow one to find the population $\rho_{0}$ (5.55) in the limit $t \rightarrow \infty$, which reads

$$
\rho_{0}(t \rightarrow \infty)= \begin{cases}e^{-2 \gamma_{0} t} e^{-2 \pi V^{2} / \alpha} ; & g<0  \tag{5.57}\\ e^{-2 \gamma_{1} t}\left(\frac{V^{2}}{\alpha^{2} t^{2}}\right)^{2} e^{2 \pi V^{2} / \alpha} ; & g>0\end{cases}
$$

One sees that apart from the exponential factor allowing for the decay of the state $|0\rangle$, the asymptotic population for $g<0$ coincides with that of (5.44) which is typical of non-decaying levels. On the contrary, if the decay rate of the state $|0\rangle$ exceeds that of the state $|1\rangle$, that is $g<0$, the asymptotic behavior is completely different, since the population $\rho_{0}$ exists mainly because of $\rho_{1}$, which decays more slowly and feeds $\rho_{0}$ via the interaction $V$ with a rate proportional to the power of the ratio of the coupling to the instantaneous detuning.


Fig. 5.4. Landau-Zener transition between two decaying levels. (a) The population of the state $|0\rangle$ for the case $\gamma_{0}<\gamma_{1}$ (normalized to the population of the decaying state at zero coupling $V=0$ ). The asymptotic value of the population is the same, as for the case of non-decaying states, although the time evolution is somewhat different. (b) The population of the state $|1\rangle$ for the case of $\gamma_{0}>\gamma_{1}$ manifests completely different behavior and exceeds sometimes considerably $\rho_{1}$ for the case of non-decaying levels, although it vanishes in the long-time asymptotic (5.58).

### 5.2.3 Landau-Zener Transition in the Presence of Transversal Relaxation

The decay of quantum state populations resulting from the interaction with a continuum or continua is just one of many examples of irreversible processes that may occur in quantum systems. Another type of relaxation takes place when a system is subjected to a perturbation which we don't know in all details, such as white noise for instance. This makes the system evolve in Hilbert space, along an unknown and unpredictable trajectory. This process can no longer be considered in terms of the Hamiltonian dynamics of a quantum state, but requires the density matrix formalism incorporating a statistical model for this unknown action. The description is usually given in terms of the density matrix $\widehat{\rho}(t)=\overline{|\Psi\rangle\langle\Psi|}$ of the system, which is the direct product of the state vector by its conjugate averaged over all possible realizations of the unknown external action. The simplest model takes into account the relaxation by introducing the so-called relaxation operator $\widehat{\hat{\mathcal{R}}}$ acting on $\widehat{\rho}(t)$, such that the master equation for the density matrix reads

$$
\begin{equation*}
i \dot{\hat{\rho}}(t)=[\widehat{H}, \widehat{\rho}(t)]+\widehat{\hat{\mathcal{R}}} \widehat{\rho}(t) \tag{5.58}
\end{equation*}
$$

as we have already seen on p. 86. The action of the relaxation operator to the $n \times n$ density matrix does not have the form of multiplication by another matrix of the same size, but it is a convolution with the forth rank $n \times n \times n \times n$ tensor $\widehat{\hat{\mathcal{R}}}$, that is $(\widehat{\hat{\mathcal{R}}} \widehat{\rho}(t))_{i j}=\sum_{k l} R_{i j k l} \rho_{k l}$. For the two-level system with the Landau-Zener Hamiltonian (5.22)

$$
\widehat{H}=\left(\begin{array}{cc}
\alpha t / 2 & V  \tag{5.59}\\
V & -\alpha t / 2
\end{array}\right)
$$

(5.58) has the form

$$
\begin{align*}
i \dot{\rho}_{00} & =V\left(\rho_{10}-\rho_{01}\right)+i \rho_{11} / T_{1} \\
i \dot{\rho}_{01} & =\alpha t \rho_{10}+V\left(\rho_{11}-\rho_{00}\right)-i \rho_{10} / T_{2} \\
i \dot{\rho}_{10} & =-\alpha t \rho_{10}+V\left(\rho_{00}-\rho_{11}\right)-i \rho_{10} / T_{2} \\
i \dot{\rho}_{11} & =V\left(\rho_{01}-\rho_{10}\right)-i \rho_{11} / T_{1} \tag{5.60}
\end{align*}
$$

where $T_{1}$ and $T_{2}$ are two parameters of the relaxation matrix called the longitudinal and transversal relaxation times.

By adding the first and the fourth equations in the system (5.60) one immediately finds the integral of motion $\rho_{11}+\rho_{00}=$ const, and the constant has to be unity as it follows from the normalization condition. The remaining three linearly independent equations rewritten in terms of three new variables: the polarization $\mathcal{P}=\rho_{01}+\rho_{10}$, dispersion $\mathcal{Q}=\rho_{01}-\rho_{10}$, and population of the upper level $\rho_{11} \equiv \frac{1}{2} \mathcal{Z}$, read

$$
\begin{align*}
& \dot{\mathcal{Z}}=2 V \mathcal{P}-\gamma \mathcal{Z} \\
& \dot{\mathcal{P}}=-2 V \mathcal{Z}+\alpha t \mathcal{Q}-\gamma \mathcal{P}+2 V \\
& \dot{\mathcal{Q}}=-\alpha t \mathcal{P}-\gamma \mathcal{Q} \tag{5.61}
\end{align*}
$$

Here we have assumed that the constant $\gamma=1 / T_{1,2}$ is identical for the transverse and longitudinal relaxation.

The initial conditions corresponding to the case $\rho_{00}(-\infty)=1$ are

$$
\begin{equation*}
\mathcal{Z}(-\infty)=\mathcal{P}(-\infty)=\mathcal{Q}(-\infty)=0 \tag{5.62}
\end{equation*}
$$

For large positive time, the two-level system is out of resonance, and hence in the asymptotic $t \rightarrow \infty$ due to the longitudinal relaxation, the system returns to its initial state $|0\rangle$. This gives the boundary condition

$$
\begin{equation*}
\mathcal{Z}(\infty)=\mathcal{P}(\infty)=\mathcal{Q}(\infty)=0 \tag{5.63}
\end{equation*}
$$

Since the asymptotic population distribution between the states is evident, we concentrate here on a non-trivial characteristic of the process - the longitudinal relaxation yield $I$, namely on the mean number of the elementary acts of the longitudinal relaxation given by the expression

$$
\begin{equation*}
I=\gamma \int_{-\infty}^{\infty} \rho_{11}(t) d t=\frac{\gamma}{2} \int_{-\infty}^{\infty} \mathcal{Z}(t) d t \tag{5.64}
\end{equation*}
$$

Equations (5.62)-(5.63) allow us to perform a Fourier transformation of (5.61) which yields

$$
\begin{align*}
i \omega \mathcal{Z} & =2 V \mathcal{P}-\gamma \mathcal{Z} \\
i \omega \mathcal{P} & =-2 V \mathcal{Z}+i \alpha \frac{\partial}{\partial \omega} \mathcal{Q}-\gamma \mathcal{P}+4 \pi V \delta(\omega) \\
i \omega \mathcal{Q} & =-i \alpha \frac{\partial}{\partial \omega} \mathcal{P}-\gamma \mathcal{Q} \tag{5.65}
\end{align*}
$$

Now we express $\mathcal{Z}$ and $\mathcal{Q}$ in terms of $\mathcal{P}$ with the help of the first and the third equations of (5.65) and substitute them into the second equation of this set. This yields

$$
\begin{equation*}
\frac{\alpha^{2}}{\omega-i \gamma} \frac{\partial}{\partial \omega} \frac{1}{\omega-i \gamma} \frac{\partial}{\partial \omega} \mathcal{P}+\mathcal{P}-\frac{4 V^{2}}{(\omega-i \gamma)^{2}} \mathcal{P}=4 \pi \frac{V}{\gamma} \delta(\omega) . \tag{5.66}
\end{equation*}
$$

The Dirac $\delta$-function on the right-hand side of the equation implies that at the point $\omega=0$ the Fourier transform of the polarization $\mathcal{P}(\omega)$ has a discontinuity in the first derivative. When we integrate (5.66) over the small vicinity $\epsilon$ of $\omega=0$ we find

$$
\begin{equation*}
\left.\frac{\partial \mathcal{P}}{\partial \omega}\right|_{\omega=0+\epsilon}-\left.\frac{\partial \mathcal{P}}{\partial \omega}\right|_{\omega=0-\epsilon}=-4 \pi \frac{\gamma V}{\alpha^{2}} . \tag{5.67}
\end{equation*}
$$

Thus, we have to find the solution of the homogeneous equation (5.66) that satisfies the boundary condition, (5.67), and corresponds to a real $\mathcal{P}(t)$, that is,

$$
\begin{equation*}
\mathcal{P}(\omega)=\mathcal{P}^{*}(-\omega) . \tag{5.68}
\end{equation*}
$$

We may now obtain the exact analytical solution of (5.66) in terms of Whittaker functions, which are another form of degenerate hypergeometric functions of (5.39).

$$
\begin{align*}
W_{\lambda, \mu}(z)= & \frac{\Gamma(-2 \mu) z^{\frac{1}{2}+\mu} e^{-\frac{z}{2}}}{\Gamma\left(\frac{1}{2}-\mu-\lambda\right)}{ }_{1} F_{1}\left(\frac{1}{2}+\mu-\lambda ; 1+2 \mu ; z\right) \\
& +\frac{\Gamma(2 \mu) z^{\frac{1}{2}-\mu} e^{-\frac{z}{2}}}{\Gamma\left(\frac{1}{2}+\mu-\lambda\right)}{ }_{1} F_{1}\left(\frac{1}{2}-\mu-\lambda ; 1-2 \mu ; z\right), \tag{5.69}
\end{align*}
$$

and have an integral representation

$$
\begin{equation*}
W_{\lambda, \mu}(z)=\frac{z^{\frac{1}{2}+\mu} e^{-\frac{z}{2}}\left(1+e^{2 i \pi(\mu-\lambda)}\right)^{-1}}{\Gamma\left(\frac{1}{2}+\mu-\lambda\right)} \int_{C} \tau^{\mu-\frac{1}{2}-\lambda}(\tau+1)^{\mu-\frac{1}{2}+\lambda} e^{-\tau x} d \tau \tag{5.70}
\end{equation*}
$$

with the integration contour $C$ going from $\infty$ to $\infty$ around the point 0 .
In order to simplify the notation, we put $2 V=1, \gamma=\widetilde{\gamma} /(2 V), \alpha=$ $\widetilde{\alpha} /(2 V)^{2}, \omega=\widetilde{\omega} /(2 V)$, and introduce the new variable $x=(i \omega+\gamma)^{2} / 2 \alpha$. Then we have

$$
\begin{equation*}
x \frac{\partial^{2}}{\partial x^{2}} \mathcal{P}+\frac{1}{2 \widetilde{\alpha}} \mathcal{P}+x \mathcal{P}=0 \tag{5.71}
\end{equation*}
$$

Solutions of the differential equations with the coefficients that are first-order polynomials in the independent variable, can be found with the help of the socalled Laplace contour integral method, which suggests that the differential equation

$$
\begin{equation*}
\sum_{n}\left(a_{n}+b_{n} x\right) \frac{\partial^{n}}{\partial x^{n}} \mathcal{Y}=0 \tag{5.72}
\end{equation*}
$$

has the solution

$$
\begin{equation*}
\mathcal{Y}(x)=\int_{C} \frac{d \tau}{\sum_{n} b_{n} \tau^{n}} \exp \left[x \tau+\int^{\tau} \frac{\sum_{n} a_{n} y^{n}}{\sum_{n} b_{n} y^{n}} d y\right] \tag{5.73}
\end{equation*}
$$

where the integration contour $C$ has to be chosen such that the exponent in the integrand has the same values at each end. Different contours satisfying this condition usually yield different, linearly independent solutions, provided they circumvent differently the singularities of the function $\left(\sum_{n} b_{n} \tau^{n}\right)^{-1}$. In particular for (5.71) one has $\left(\sum_{n} b_{n} \tau^{n}\right)^{-1}=\left(\tau^{2}+1\right)^{-1}$; $\int^{\tau} \sum_{n} a_{n} y^{n} / \sum_{n} b_{n} y^{n} d y=\frac{1}{4 \widetilde{\alpha}} \ln \left(\tau^{2}+1\right)$, and this method yields

$$
\begin{equation*}
\mathcal{P}=\frac{A^{\prime}}{2 i \alpha} \int_{C_{W}} \frac{\exp (\tau x)}{(\tau+i)^{1-i / 4 \widetilde{\alpha}}(\tau-i)^{1+i / 4 \widetilde{\alpha}}} d \tau \tag{5.74}
\end{equation*}
$$

The contour $C_{W}$ arrives from $\tau=i \infty$, goes around the point $\tau=i$, returning along the direction $\tau=i \infty$. This integration path ensures convergence of the integral for $\omega>0$. The integral representation for $\omega<0$ follows from the condition (5.68). We note that another possible contour $C_{M}$ is the loop which circumvents both points $i$ and $-i$. The corresponding polarization $\mathcal{P}$ accounts for the homogeneous solution which is regular at the point $\omega=0$ and therefore does not contribute to the final result.

The substitution $\tau=i+i 2 z$ transforms the integral (5.74) to the integral representation of the Whittaker function (5.70),

$$
\begin{equation*}
u_{1}(\omega)=W_{i / 4 \alpha,-1 / 2}\left(-i \frac{(i \omega+\gamma)^{2}}{\alpha}\right) \tag{5.75}
\end{equation*}
$$

Here, we have made use of the explicit expression for $x$. Hence, the Whittaker function $u_{1}(\omega)$ is a solution of (5.71) for $\omega>0$, that is, $\left.\mathcal{P}(\omega)\right|_{\omega>0}=A u_{1}(\omega)$. The other linearly independent solution of (5.66),

$$
\begin{equation*}
\mathcal{P}(\omega)=A^{*} u_{2}(\omega)=A^{*} W_{-i / 4 \alpha,-1 / 2}\left(i \frac{(i \omega+\gamma)^{2}}{\alpha}\right) \tag{5.76}
\end{equation*}
$$

accounts for the case $\omega<0$.
Now we can find the proportionality constant $A$ with the help of the condition given in (5.68) at $\omega=0$ and the condition (5.67). We come then to the set of equations

$$
\begin{align*}
& A u_{1}(0)-A^{*} u_{2}(0)=0 \\
& \left.A \frac{\partial u_{1}}{\partial \omega}\right|_{\omega=0}-\left.A^{*} \frac{\partial u_{2}}{\partial \omega}\right|_{\omega=0}=-2 \pi \frac{\gamma}{\alpha^{2}} . \tag{5.77}
\end{align*}
$$

The determinant $D$ of this set is proportional to the Wronskian $W\left(u_{1} ; u_{2}\right) \equiv$ $\left(u_{1} \frac{\partial}{\partial x} u_{2}-u_{2} \frac{\partial}{\partial x} u_{1}\right)$ of the Whittaker functions and reads

$$
\begin{equation*}
D=\frac{2 \gamma}{\alpha} W\left(u_{1} ; u_{2}\right)=\frac{2 \gamma}{\alpha} \exp \left(\frac{\pi}{4 \alpha}\right) . \tag{5.78}
\end{equation*}
$$

The factor $\frac{2 \gamma}{\alpha}$ results from the derivative of the arguments $\frac{\partial x}{\partial \omega}$. Solving (5.77) we obtain the constant $A$ and find for the polarization

$$
\begin{equation*}
\left.\mathcal{P}(\omega)\right|_{\omega>0}=\pi \frac{u_{2}(0) u_{1}(\omega)}{\alpha} \exp \left(-\frac{\pi}{4 \alpha}\right) . \tag{5.79}
\end{equation*}
$$

Now we are able to calculate the longitudinal relaxation yield (5.64) that is $I=\frac{\gamma}{2} \int_{-\infty}^{\infty} \mathcal{Z}(t) d t=\frac{\gamma}{2} \mathcal{Z}(\omega=0)$. Since, according to the first line of (5.65), we have $\frac{1}{2} \mathcal{Z}(\omega=0)=(2 \gamma)^{-1} \mathcal{P}(\omega=0)$ we find from Eqs.(5.75)-(5.76) and (5.79) for $\omega=0$

$$
\begin{equation*}
I=2 \pi \frac{V^{2}}{\alpha} e^{\frac{-\pi V^{2}}{\alpha}}\left|W_{\frac{i V^{2}}{\alpha},-\frac{1}{2}}\left(-\frac{i \gamma^{2}}{\alpha}\right)\right|^{2} . \tag{5.80}
\end{equation*}
$$

In Fig. 5.5 we show the dependence of the longitudinal relaxation yield $I$ on the adiabaticity parameter $V^{2} / \alpha$ and on the quenching parameter $\gamma^{2} / \alpha$ suggested by (5.80). At $\gamma=0$, the yield is given by the regular Landau-Zener dependence (5.44) for $\rho_{1}=1-\rho_{0}$, which simply means that the relaxation yield amounts to the probability to make the transition. It changes with the increase of the quenching parameter $\gamma^{2} / \alpha$, as the relaxation yield increases with increasing relaxation rate.

We now consider the asymptote of (5.80). We start with the case of slow passage through the resonance, that is $\alpha \rightarrow 0$. We use the integral representation of the Whittaker function and find the following expression with the help of the stationary phase method

$$
\begin{align*}
W_{\frac{i V^{2}}{4 \alpha},-\frac{1}{2}}\left(-\frac{i \gamma^{2}}{\alpha}\right) & =\frac{e^{i \gamma^{2} / 2 \alpha}}{\Gamma\left(-\frac{i V^{2}}{\alpha}\right)} \int_{0}^{\infty} \frac{\exp \left\{i \frac{\gamma^{2}}{\alpha} t-i \frac{V^{2}}{\alpha} \ln \frac{t}{1+t}\right\}}{t(t+1)} d t \\
& \simeq \frac{e^{i \Phi}}{\Gamma\left(-i V^{2} / \alpha\right)} \frac{\sqrt{\pi \alpha}}{\left.\left[1+(2 V / \gamma)^{2}\right)\right]^{1 / 4} V} \tag{5.81}
\end{align*}
$$

where $\Phi$ is a phase factor that does not effect the final result, (5.80). Here we have taken into account that in the limit $\alpha \rightarrow 0$ the only stationary point on the integration path is $t=-\frac{1}{2}+\frac{1}{2} \sqrt{1+\frac{4 V^{2}}{\gamma^{2}}}$. We also have assumed $V^{2} \gamma^{2} / \alpha^{2} \gg 1$, allowing us to neglect derivatives higher than second order at


Fig. 5.5. Landau-Zener transition to a decaying level. (a) Level scheme. (b) Expectation value of the relaxation yield $I$ given by (5.80) as a function of the LandauZener parameter $V^{2} / \alpha$ and quenching parameter $\gamma^{2} / \alpha$. At $\gamma=0$ the dependency recovers the probability of the Landau-Zener transition $\rho_{1}=1-\rho_{0}$.
the stationary phase point and thus satisfy the requirement of the stationary phase method. Substitution of (5.81) into (5.80) results in

$$
\begin{equation*}
I=2 \pi \frac{\gamma V^{2}}{\alpha} \frac{1-\exp \left\{\frac{-2 \pi V^{2}}{\alpha}\right\}}{\left(\gamma^{2}+4 V^{2}\right)^{1 / 2}} \tag{5.82}
\end{equation*}
$$

where we have made use of (5.46). The exponential term, reminiscent of the Landau-Zener expression, arises from the $\Gamma$-function and not from the stationary phase analysis of (5.81). However, in order to be consistent within the order of approximation it should be neglected. Note that this result also follows from the rate equation. Indeed, one can obtain the result (5.82) directly from (5.61) in the quasi-stationary limit, namely when all time derivatives are equal to zero. This results in $\mathcal{Z}=2 V \mathcal{P} / \gamma$, and $\mathcal{Q}=-\alpha t \mathcal{P} / \gamma$, which gives $\mathcal{P}=2 \gamma V\left[4 V^{2}+(\alpha t)^{2}+\gamma^{2}\right]^{-1}$ and hence

$$
I=V \int \mathcal{P} d t=2 \pi \frac{\gamma V^{2}}{\alpha \sqrt{\gamma^{2}+4 V^{2}}}
$$

The limit of weak interaction, $V^{2} / \alpha \rightarrow 0$, follows from the condition $\left|W_{0,-1 / 2}(x)\right|=\left|e^{i x}\right|=1$, and (5.80) simplifies to $I \simeq 2 \pi V^{2} / \alpha$. This implies that the asymptotic formula (5.82) is also valid for weak interactions.

We now consider the limit of slow decay, $\gamma^{2} / \alpha \ll 1$, with the help of the asymptotic expression for the Whittaker function

$$
\begin{equation*}
\lim _{z \rightarrow 0} W_{i a,-\frac{1}{2}}(z)=\frac{1}{\Gamma(1+i a)}=\frac{1}{i a \Gamma(i a)} \tag{5.83}
\end{equation*}
$$

which, after substitution into (5.80), results in

$$
\begin{equation*}
I=1-\exp \left(-\frac{2 \pi V^{2}}{\alpha}\right) \tag{5.84}
\end{equation*}
$$

We hence recover the probability of the Landau-Zener transition, (5.44), which indeed is identical to the probability of emission of a spontaneous photon in the case of a small decay rate. We note that a simple expression for the limit $V^{2} \gamma^{2} / \alpha^{2} \ll 1$ is difficult to obtain due to the non-analytical behavior of the Whittaker function at small values of the argument, as shown in Fig. 5.2.

We conclude by noting that Rabi oscillations of population between the states (Fig. 5.2) that take place while the levels are in resonance, do not manifest themselves in any oscillations of the asymptotic populations nor in the relaxation yield. At first sight, one might think that the relaxation would bring them to light, since it "brakes" the population oscillations in the middle of the Rabi cycle, that is, before the passage through the resonance is complete. However, they remain hidden even in the case when relaxation processes are taken into account: the irreversible decay does not affect the probability of the transition at all as long as the population of the initially populated state decays faster then the population of the other state. The longitudinal relaxation of the upper level accounts for the multiple returns of the particle to the initial state in the middle of the Rabi cycle. However, it is also the case that these returns do not bring any oscillations to the expectation value of the relaxation yield. The gradual change of the quenching parameter results instead in the gradual transformation of the Landau-Zener dependence to the smooth dependence given by the rate equation.

### 5.2.4 Excitation by a Pulse

Another example of an exactly soluble two-level problem corresponds to the case of a constant detuning $\Delta_{0}=0, \Delta_{1}=$ const and the time-dependent interaction of a very special profile $V(t)=V / \cosh \alpha t$. The Schrödinger equation (5.29) after the replacements

$$
\begin{equation*}
z=\frac{1}{1+e^{-2 \alpha t}} ; \quad \Delta=\frac{\Delta_{1}}{2 \alpha} ; \quad \theta=\frac{V \pi}{\alpha} \tag{5.85}
\end{equation*}
$$

yields the second-order differential equation (5.30) in the form

$$
\begin{equation*}
-z(z-1) \frac{\partial^{2} \psi_{0}}{\partial z^{2}}+\left(\frac{1}{2}-z+i \Delta\right) \frac{\partial \psi_{0}}{\partial z}+\frac{2 \theta}{\pi} \psi_{0}=0 \tag{5.86}
\end{equation*}
$$

which for $\psi_{0}(t \rightarrow-\infty)=\psi_{0}(z \rightarrow 0) \rightarrow 1$ has an explicit solution in terms of the hypergeometric function

$$
\begin{equation*}
\psi_{0}(z)={ }_{2} F_{1}\left(\frac{\theta}{\pi},-\frac{\theta}{\pi} ; \frac{1}{2}+i \Delta ; z\right) . \tag{5.87}
\end{equation*}
$$

This expression yields the amplitude as $t \rightarrow \infty$

$$
\begin{align*}
\left.\psi_{0}(t)\right|_{t \rightarrow \infty} & =\left.\psi_{0}(z)\right|_{z=1}={ }_{2} F_{1}\left(\frac{\theta}{\pi},-\frac{\theta}{\pi} ; \frac{1}{2}+i \Delta ; 1\right) \\
& =\frac{\Gamma\left(\frac{1}{2}+i \Delta\right) \Gamma\left(\frac{1}{2}+i \Delta\right)}{\Gamma\left(\frac{1}{2}+i \Delta-\frac{\theta}{\pi}\right) \Gamma\left(\frac{1}{2}+i \Delta+\frac{\theta}{\pi}\right)} \tag{5.88}
\end{align*}
$$

in terms of $\Gamma$-functions. The corresponding population reads

$$
\begin{equation*}
\left.\rho_{0}(t)\right|_{t \rightarrow \infty}=\frac{\cosh (2 \pi \Delta)+\cos (2 \theta)}{2 \cosh ^{2}(\pi \Delta)} \tag{5.89}
\end{equation*}
$$

where (5.46) and the relations

$$
\begin{align*}
\Gamma(x+1) & =x \Gamma(x) \\
|\Gamma(1-x+i \Delta) \Gamma(1+x+i \Delta)|^{2} & =\frac{2 \pi^{2}\left(x^{2}+\Delta^{2}\right)}{\cosh (2 \pi \Delta)-\cos (2 \pi x)} \tag{5.90}
\end{align*}
$$

have been employed. For the transition probability, after trigonometric transformations one obtains

$$
\begin{equation*}
\left.\rho_{1}(t)\right|_{t \rightarrow \infty}=\frac{\sin ^{2} \theta}{\cosh ^{2}(\pi \Delta)} \tag{5.91}
\end{equation*}
$$

In Fig. 5.6 we show the time-dependent transition probability $1-\left|\psi_{0}(t)\right|^{2}$ corresponding to the amplitude (5.87), this amplitude, and the asymptotic


Fig. 5.6. Excitation of a two-level system by a pulse $V / \cosh ^{-2}(2 \alpha t)$. (a) Transition probability as a function of time $T=2 \alpha t$ and detuning $\Delta=\Delta_{1} / 2 \alpha$. One sees oscillations in resonance, and the adiabatic regime for large detunings. (b) The asymptotic probability as a function of the detuning and the pulse area $\theta=V / \alpha$. (c) Real (dashed line) and imaginary (dot-dash line) parts of the probability amplitude to remain in the state $|0\rangle$ as a function of time for fixed detuning $\Delta_{1} / 2 \alpha=0.7$ and pulse area $\theta=2.45$.
probability transition (5.91).
We note one very specific property of the problem under consideration. Unlike other cases of time-dependence, for the interaction profile $V / \cosh \alpha t$ the transition probability is a product of two factors, one of which is a function of the detuning $\Delta$ only, whereas the other depends only on the coupling amplitude $V$. Moreover, the last enters the expression in the form of the socalled pulse area $\theta=\int_{-\infty}^{\infty} V / \cosh (\alpha t) d t$. Each of these two factors has clear physical meaning. Indeed, for $\Delta_{0 ; 1}=0$ equations (5.29) take a particularly simple form $i \dot{\psi}_{0}=V(t) \psi_{1}, i \dot{\psi}_{1}=V(t) \psi_{0}$, and have an exact solution

$$
\begin{align*}
& \psi_{0}(t)=\cos \left[\int_{-\infty}^{t} V\left(t^{\prime}\right) d t^{\prime}\right] \\
& \psi_{1}(t)=i \sin \left[\int_{-\infty}^{t} V\left(t^{\prime}\right) d t^{\prime}\right] \tag{5.92}
\end{align*}
$$

which yields $\psi_{1}(\infty)=i \sin \theta$, and hence the expression $\rho_{1}(\infty)=\sin ^{2} \theta$ coinciding completely with (5.91) in resonance. The other factor allows for non-zero detunings, being just the line shape $\left|V_{\omega}\right|^{2}$ at $\omega=\Delta_{1}$ given by the Fourier transform of the $V / \cosh \alpha t$ dependence and normalized by the condition $V_{\omega=0}=1$. Indeed, employing the Cauchy formula for residuals at the points $\alpha t=i(\pi / 2+n \pi)$, for $n=0,1 \ldots$, yields

$$
\begin{align*}
\frac{1}{\pi i} \int_{-\infty}^{\infty} \frac{\exp \left(i \Delta_{1} t\right)}{\cosh (\alpha t)} d t & =2 \sum_{n=0}^{\infty}(-1)^{n} \exp \left(-\frac{\Delta_{1}}{\alpha}(\pi / 2+n \pi)\right) \\
& =\frac{2 \exp \left(-\frac{\Delta_{1} \pi}{2 \alpha}\right)}{1+\exp \left(-\frac{\Delta_{1} \pi}{\alpha}\right)}=\cosh ^{-1}\left(\frac{\Delta_{1} \pi}{2 \alpha}\right) \tag{5.93}
\end{align*}
$$

We emphasize that such a simple asymptotic form of (5.91) is more a coincidence than the result of separation between the effects of detuning and coupling that remain strongly entangled during the evolution at finite $t$. One can see this in Fig.5.6(a), where the difference of the time behavior of the population $\rho_{1}(t)$ for different detunings $\Delta$ is clearly seen.

### 5.2.5 Exponentially Rising Coupling

We consider one more example of exactly soluble problems for which the interaction $V(t)$ between the states rises exponentially, while the detuning $\Delta_{1}$ remains constant, also known as the Demkov-Olson problem, for which the Schrödinger equation (5.29) reads

$$
\begin{align*}
i \dot{\psi}_{0} & =V e^{\alpha t} \psi_{1} \\
i \dot{\psi}_{1} & =\Delta_{1} \psi_{1}+V e^{\alpha t} \psi_{0} \tag{5.94}
\end{align*}
$$

We set $\Delta_{0}=0$ as earlier, and introduce the variable $\theta=\int^{t} V e^{\alpha x} d x=V e^{\alpha t} / \alpha$ such, that the system takes the form

$$
\begin{align*}
& i \frac{\partial \psi_{0}}{\partial \theta}=\psi_{1} \\
& i \frac{\partial \psi_{1}}{\partial \theta}=\frac{\Delta_{1}}{\alpha \theta} \psi_{1}+\psi_{0} \tag{5.95}
\end{align*}
$$

where the initial condition is $\psi_{0}(\theta=0)=1$.
Now one can either employ the transformation (5.30) that yields the Bessel equation, or look for the solution in the form of a power series $\psi_{0}(\theta)=\sum_{n=0}^{\infty} c_{2 n} \theta^{2 n}$ and $\psi_{1}(\theta)=\sum_{n=0}^{\infty} c_{2 n+1} \theta^{2 n+1}$, as has been done for (5.37). The last method immediately allows one to incorporate the initial condition by setting $c_{0}=1$, and after making use of the recurrence relations

$$
\begin{align*}
i 2 n c_{2 n} & =c_{2 n-1}, \\
i(2 n-1) c_{2 n-1} & =\frac{\Delta_{1}}{\alpha} c_{2 n-1}+c_{2(n-1)}, \tag{5.96}
\end{align*}
$$

that is $c_{2 n}=-c_{2(n-1)} / 2 n\left(2 n-1+i \frac{\Delta_{1}}{\alpha}\right)$, this approach results in

$$
\begin{equation*}
c_{2 n}=\left(\frac{-1}{4}\right)^{n} \frac{1}{n!} \frac{\Gamma\left(\frac{1}{2}+\frac{i \Delta_{1}}{2 \alpha}\right)}{\Gamma\left(n+\frac{1}{2}+\frac{i \Delta_{1}}{2 \alpha}\right)} . \tag{5.97}
\end{equation*}
$$

This means that

$$
\begin{align*}
\psi_{0}(\theta) & =\Gamma\left(\frac{1}{2}+\frac{i \Delta_{1}}{2 \alpha}\right) \sum_{n=0}^{\infty} \frac{1}{n!} \frac{\left(-\theta^{2} / 4\right)^{n}}{\Gamma\left(n+\frac{1}{2}+\frac{i \Delta_{1}}{2 \alpha}\right)} \\
& =\Gamma\left(\frac{1}{2}+i \Delta\right)\left(\frac{\theta}{2}\right)^{1 / 2-i \Delta} J_{i \Delta-1 / 2}(\theta) \tag{5.98}
\end{align*}
$$

where a standard representation of the Bessel function $J_{\nu}(\theta)$ in the form of a power series has been employed, and where we make use of the notation (5.85), although now the parameter $\alpha$ corresponds to an exponential, and not the $1 / \cosh x$ dependence, and even has a different dimensionality. From first equation of the set (5.95) and with allowance of the rule for the Bessel function derivatives one also obtains

$$
\begin{equation*}
\psi_{1}(\theta)=-i \Gamma\left(\frac{1}{2}+i \Delta\right)\left(\frac{\theta}{2}\right)^{1 / 2-i \Delta} J_{i \Delta+1 / 2}(\theta) \tag{5.99}
\end{equation*}
$$

Equation (5.98) with allowance of (5.46) gives the probabilities

$$
\begin{align*}
\rho_{0}(\theta) & =\left|\Gamma\left(\frac{1}{2}+i \Delta\right)\right|^{2} \frac{\theta}{2}\left|J_{i \Delta-1 / 2}(\theta)\right|^{2} \\
& =\frac{\pi \theta}{2} \cosh ^{-1}(\pi \Delta)\left|J_{i \Delta-1 / 2}(\theta)\right|^{2}, \\
\rho_{1}(\theta) & =\frac{\pi \theta}{2} \cosh ^{-1}(\pi \Delta)\left|J_{i \Delta+1 / 2}(\theta)\right|^{2} . \tag{5.100}
\end{align*}
$$



Fig. 5.7. Excitation of a two-level system by an exponentially increasing interaction $V / e^{\alpha t}$. (a) Transition probability as a function of the time parameter $\theta=V e^{\alpha t} / \alpha$ and the detuning $\Delta=\Delta_{1} / 2 \alpha$. One sees oscillations for small detunings and the stationary regime with $\rho_{0 ; 1} \rightarrow 1 / 2$. (b) The probability of the adiabatic transition as a function of $T=\alpha t$ and $\ln \left(\Delta_{1} / 2 \alpha\right)$. At small detunings one can identify a regime where the non-adiabatic transition probability tends to $1 / 4$.

The population of the state $|0\rangle$ is shown in Fig. 5.7. It is also interesting to trace the populations of the adiabatic states, that is the instantaneous eigenstates $| \pm\rangle$ of the Hamiltonian, that can be found by analogy to (5.48) as

Let us assume $\Delta>0$ and denote by $\Theta$ the angle in Hilbert space between the state vector $|+\rangle$ and the state $|0\rangle$, that is

$$
\cos \Theta=\frac{\left(\Delta+\sqrt{\Delta^{2}+\theta^{2}}\right)}{\sqrt{\left(\Delta+\sqrt{\Delta^{2}+\theta^{2}}\right)^{2}+\theta^{2}}}
$$

and find the population

$$
\begin{equation*}
\rho_{-}(\theta)=\frac{\pi \theta\left|J_{i \Delta-1 / 2}(\theta) \cos \Theta-i J_{i \Delta+1 / 2}(\theta) \sin \Theta\right|^{2}}{2 \cosh (\pi \Delta)} \tag{5.102}
\end{equation*}
$$

of the state $|-\rangle$ that determines the probability of a non-adiabatic transition. In Fig. 5.7(b) we depict this population as a function of time and the detuning. One sees that for small detunings in the long-time regime it approaches the asymptotic value $1 / 4$, which can be seen directly from (5.102) with the help of the asymptotic expression

$$
\begin{equation*}
J_{i \Delta \pm 1 / 2}(\theta) \simeq \sqrt{\frac{1}{\pi \theta}} \cosh \left(i \theta+\Delta \pi-(1 \pm 1) \frac{\pi i}{4}\right) \tag{5.103}
\end{equation*}
$$

and with allowance of the fact that $\cos \Theta \simeq \sin \Theta \simeq 1 / \sqrt{2}$ as $t \rightarrow \infty$ for large $\theta$ and interactions. The substitution yields

$$
\begin{equation*}
\rho_{n a}(\theta \rightarrow \infty)=\frac{\exp (-2|\Delta| \pi)}{4 \cosh (\pi \Delta)} \tag{5.104}
\end{equation*}
$$

where we have also taken into account that for negative $\Delta$ the probability of the non adiabatic transition $\rho_{n a}(\theta)$ is given by $\rho_{+}$.

We note one general property common for all two-level problems. For slow variation of the system $\alpha \rightarrow 0$ and the probability to change the adiabatic state given by (5.50), (5.91), and (5.104) is always small.

### 5.3 Semiclassical Analysis of Time-Dependent Systems

Considering the Landau-Zener problem on p. 198, we have seen that the populations of adiabatic states change slightly when the Hamiltonian varies slowly. This is a general situation, also valid in the case of a large system consisting of many quantum levels. For the description of the population dynamics, one may make use of a treatment similar to the semiclassical approach to the Schrödinger equation of a particle moving in an external potential, with the only difference that time $t$ plays the role of the coordinate. A straightforward substitution of the semiclassical ansatz $\exp [-i S(t)]$ either into (5.36) or (5.37) results however in a complex-valued solution for $S(t)$, which means that there is no direct classical analog for the mechanical action in these representations. Moreover, the Hamilton-Jacobi equation, obtained for the classical action $S(t)$ after the substitution $\exp [-i S(t)]$ into the Schrödinger equation, turns out to be dependent on the choice of representation, while the allowance for the second derivative of $S(t)$ via perturbation theory does not improve the situation much. Therefore, in order to make use of the semiclassical analysis, one needs first to find a representation where the equation for the action would have an appropriate form that would demonstrate the vanishing transition probabilities among the adiabatic states.

### 5.3.1 Two-Level Systems and the Dykhne Formula

We start the analysis with the semiclassical consideration of the LandauZener problem and determine the reference system most suitable for the semiclassical consideration. To this end we employ the adiabatic representation by finding a time-dependent transformation $\widehat{O}(t)$ diagonalizing the Hamiltonian (5.32) by the transformation $\widehat{H^{\prime}}(t)=\widehat{O}^{-1}(t) \widehat{H}(t) \widehat{O}(t)$, which corresponds to the rotation from our initial diabatic representation to the adiabatic basis.

For the Landau-Zener problem, the unitary matrix performing the rotation has the form

$$
\begin{equation*}
\widehat{O}(t)=\exp \left[i \widehat{\sigma}_{y} \Theta(t) / 2\right] \tag{5.105}
\end{equation*}
$$

where $\Theta(t)=\pi / 2+\arctan (\alpha t / 2 V)$ is the angle of rotation around the $y$ axis that aligns the vector ( $V, 0, \alpha t / 2$ ) with the $z$-axis, and the Pauli matrices $\widehat{\sigma}_{x}, \widehat{\sigma}_{y}, \widehat{\sigma}_{z}$ are given explicitly by (5.33). The unitary matrix (5.105) transforms the Hamiltonian $\widehat{H}=V \widehat{\sigma}_{x}+\alpha t \widehat{\sigma}_{z} / 2$ to the diagonal form $\widehat{H}^{\prime}=\sqrt{(\alpha t / 2)^{2}+V^{2}} \hat{\sigma}_{z}$

We note here that $\widehat{O}(t)$ of (5.105) equals the unit matrix at $t=-\infty$, and at $t=\infty$ it reads

$$
\widehat{O}(\infty)=e^{i \widehat{\sigma}_{y} \pi / 2}=\left(\begin{array}{rr}
0 & 1  \tag{5.106}\\
-1 & 0
\end{array}\right),
$$

which means that not only the amplitudes of the states $|0\rangle$ and $|1\rangle$ have been interchanged after the diabatic states have passed through the resonance, but also that an additional phase factor $e^{i \pi}=-1$ appears as a result of this change in the relative position of the levels. This is a particular example of the general phenomenon known as topological or Berry phases associated with rotation of the adiabatic energy eigenstates with respect to the original basis in Hilbert space.

By substituting $|\psi\rangle=\widehat{O}(t)\left|\psi^{\prime}\right\rangle$ into (5.31)

$$
\begin{equation*}
i \frac{\partial}{\partial t} \widehat{O}(t)\left|\psi^{\prime}\right\rangle=i \frac{\widehat{\partial O(t)}}{\partial t}\left|\psi^{\prime}\right\rangle+i \widehat{O}(t) \frac{\partial}{\partial t}\left|\psi^{\prime}\right\rangle=\widehat{H}(t) \widehat{O}(t)\left|\psi^{\prime}\right\rangle \tag{5.107}
\end{equation*}
$$

one obtains the Schrödinger equation for the state vector $\left|\psi^{\prime}\right\rangle$ in the adiabatic representation

$$
\begin{equation*}
i \frac{\partial}{\partial t}\left|\psi^{\prime}\right\rangle=\left[\widehat{H}^{\prime}(t)-\widehat{C}(t)\right]\left|\psi^{\prime}\right\rangle \tag{5.108}
\end{equation*}
$$

where $\widehat{H}^{\prime}(t)=\sqrt{V^{2}+(\alpha t / 2)^{2}} \widehat{\sigma}_{z}$ is a diagonal matrix, and

$$
\widehat{C}(t)=i \widehat{O}^{-1}(t) \widehat{\dot{O}}(t)=-\widehat{\sigma}_{y} \dot{\Theta}(t) / 2
$$

In other words the Schrödinger equation has the explicit form

$$
\begin{align*}
i \frac{\partial}{\partial t} \psi_{0}^{\prime} & =-\sqrt{V^{2}+\left(\frac{\alpha t}{2}\right)^{2}} \psi_{0}^{\prime}-\frac{i V \alpha}{4 V^{2}+(\alpha t)^{2}}
\end{align*} \psi_{1}^{\prime} .
$$

One sees that in spite of the fact that the matrix $\widehat{O}(t)$ diagonalizes the Hamiltonian, the cross-terms between the adiabatic states persist due to the time dependence of this transformation, and the rotation rate $\dot{\theta}(t)$ of the adiabatic reference frame represents the probability amplitude of this process.

We now express $\psi_{1}^{\prime}$ from the first equation (5.109) and substitute it into the second one, which yields

$$
\begin{equation*}
\frac{\partial^{2} \psi_{0}^{\prime}}{\partial t^{2}}+\frac{\alpha^{2} t / 2}{V^{2}+\frac{\alpha^{2} t^{2}}{4}} \frac{\partial \psi_{0}^{\prime}}{\partial t}=\left[\frac{3 i \alpha^{2} t / 4}{\sqrt{V^{2}+\frac{\alpha^{2} t^{2}}{4}}}-V^{2}-\frac{\alpha^{2} t^{2}}{4}-\frac{\alpha^{2} / 16}{V^{2}+\frac{\alpha^{2} t^{2}}{4}}\right] \psi_{0}^{\prime} \tag{5.110}
\end{equation*}
$$

take the amplitude in the form $\psi_{0}^{\prime}=\Psi_{0}\left(4 V^{2}+(\alpha t)^{2}\right)^{-1 / 2}$ in order to get rid of the coefficient in front of the first derivative, and arrive at

$$
\begin{equation*}
\frac{\partial^{2} \Psi_{0}}{\partial t^{2}}+\left[\frac{3 i \alpha^{2} t / 2}{\sqrt{4 V^{2}+(\alpha t)^{2}}}+V^{2}+\frac{(\alpha t)^{2}}{4}+\frac{5 \alpha^{2} / 4}{\left(4 V^{2}+(\alpha t)^{2}\right)^{2}}\right] \Psi_{0}=0 \tag{5.111}
\end{equation*}
$$

The last term in square brackets can be omitted when the adiabaticity parameter $\alpha / V^{2}$ is small, and we obtain the WKB solution

$$
\begin{equation*}
\psi_{0}^{\prime}=\frac{C \exp \left[-i \int\left(\frac{3 i \alpha^{2} t / 2}{\sqrt{4 V^{2}+(\alpha t)^{2}}}+V^{2}+\frac{(\alpha t)^{2}}{4}\right)^{1 / 2} d t\right]}{\left(\frac{3 i \alpha^{2} t / 2}{\sqrt{4 V^{2}+(\alpha t)^{2}}}+V^{2}+\frac{(\alpha t)^{2}}{4}\right)^{1 / 4}\left(4 V^{2}+(\alpha t)^{2}\right)^{1 / 2}} \tag{5.112}
\end{equation*}
$$

where $C$ is a constant.
In the adiabatic condition, the imaginary term is always small and the square root in the integrand can be developed as a Taylor series, which after integration yields

$$
\begin{equation*}
\psi_{0}^{\prime}=\frac{C \exp \left[-i \int \sqrt{V^{2}+\frac{(\alpha t)^{2}}{4}} d t+\int \frac{3 \alpha^{2} t / 2}{4 V^{2}+(\alpha t)^{2}} d t\right]}{\left(V^{2}+\frac{(\alpha t)^{2}}{4}\right)^{1 / 4}\left(4 V^{2}+(\alpha t)^{2}\right)^{1 / 2}}=C^{\prime} e^{-i \int \sqrt{V^{2}+(\alpha t)^{2} / 4} d t} \tag{5.113}
\end{equation*}
$$

where to be consistent with the order of approximation we have omitted terms in the pre-exponential factors that are small and vanishing as $t \rightarrow \infty$ and have substituted the explicit expression $\int \frac{3 \alpha^{2} t / 2}{4 V^{2}+(\alpha t)^{2}} d t=(3 / 4) \ln \left[4 V^{2}+(\alpha t)^{2}\right]$. The constant $C^{\prime}$ equals unity for the initial conditions $\rho_{0}(t \rightarrow-\infty) \rightarrow 1$.

From the result (5.113) one sees that the adiabatic representation has a remarkable property - all of the first-order corrections over the parameter $\alpha / V^{2}$ within the approximation given by the diagonal Hamiltonian $\widehat{H}^{\prime}(t)$ cancel each other, such that the latter gives much better precision than is suggested by the linear perturbation theory for other representations. The reason for this is in the fact, that the probability amplitudes oscillate rapidly as compared to the typical time dependence of the non-diagonal matrix elements $i V \alpha\left(4 V^{2}+(\alpha \tau)^{2}\right)^{-1}$ responsible for the transition in (5.109). But how do we determine the accuracy of the approximation? Substituting (5.113)
into the second equation of the set (5.109) one finds the other probability amplitude $\psi_{1}^{\prime}$

$$
\begin{equation*}
\psi_{1}^{\prime}=i e^{i \Phi(t)} \int^{t} \frac{i V \alpha}{4 V^{2}+(\alpha \tau)^{2}} e^{-2 i \Phi(\tau)} d \tau \tag{5.114}
\end{equation*}
$$

Evaluation of the integrals $\Phi(t)=\int^{t} \sqrt{V^{2}+\left(\frac{\alpha t}{4}\right)^{2}} d t$ yields $\frac{t}{2} \sqrt{V^{2}+\left(\frac{\alpha t}{4}\right)^{2}}+$ $\frac{V^{2}}{\alpha} \ln \left[\frac{t \alpha}{2 V}+\sqrt{1+\left(\frac{\alpha t}{2 V}\right)^{2}}\right]$, and hence (5.114) takes the form

$$
\begin{equation*}
\psi_{1}^{\prime}=-e^{i \Phi(t)} \int^{t} \frac{V \alpha\left(\frac{\tau \alpha}{2 V}+\sqrt{1+\left(\frac{\alpha \tau}{2 V}\right)^{2}}\right)^{-i 2 V^{2} / \alpha}}{4 V^{2}+(\alpha \tau)^{2}} e^{-i \tau \sqrt{V^{2}+(\alpha \tau)^{2} / 4}} d \tau \tag{5.115}
\end{equation*}
$$

and results for $t \rightarrow \infty$, after the substitution $\tau=2 V \sinh [\chi] / \alpha$, in

$$
\begin{equation*}
\left|\psi_{1}^{\prime}\right|^{2}=\left|\int_{-\infty}^{\infty} \frac{\exp \left[\frac{V^{2}}{i \alpha}(2 \chi+\sinh [2 \chi])\right]}{2 \cosh [\chi]} d \chi\right|^{2} \tag{5.116}
\end{equation*}
$$

By analogy to (5.93), the integral in (5.116) has simple poles at the points $\chi=-\pi i / 2-i k \pi$, where $\cosh [\chi]=\sinh [2 \chi]=0$. We have to take positive $k$, corresponding to the vanishing residuals, in order to obtain the result

$$
\begin{equation*}
\left|\psi_{1}^{\prime}\right|^{2}=\left|\sum_{k=0}^{\infty} \frac{\pi \exp \left[-\frac{\pi V^{2}}{\alpha}(1+2 k)\right]}{-i(-1)^{k}}\right|^{2}=\left|\frac{i \pi \exp \left[-\frac{\pi V^{2}}{\alpha}\right]}{2\left(1+\exp \left[-\frac{2 \pi V^{2}}{\alpha}\right]\right)}\right|^{2} \simeq \pi^{2} e^{-\frac{2 \pi V^{2}}{\alpha}} \tag{5.117}
\end{equation*}
$$

coinciding with the exact result of (5.44) for the adiabatic limit $\frac{2 \pi V^{2}}{\alpha} \gg 1$ under consideration. Note that the exponentially small probability of the non-adiabatic transition has an evident physical explanation: the frequency spectrum of the time-dependent perturbation $V \alpha\left(4 V^{2}+(\alpha t)^{2}\right)^{-1}$, given by (5.109) and originating from the slow rotation of the adiabatic basis, simply has vanishing intensities for the harmonics corresponding to the transition energies between the adiabatic states, even at the moment of their closest approach.

This result can be generalized to an arbitrary time-dependent two-level Hamiltonian, provided the typical width of the frequency spectrum associated with the rotation of the adiabatic frame remains much smaller than the transition frequency between the adiabatic levels. For the Hamiltonian

$$
\begin{equation*}
\widehat{H}=H_{z}(t) \widehat{\sigma}_{z}+H_{x}(t) \widehat{\sigma}_{x} \tag{5.118}
\end{equation*}
$$

the equation (5.114) takes the form

$$
\begin{equation*}
\psi_{1}^{\prime}=e^{i \Phi(t)} \int^{t} \frac{\dot{\Theta}(\tau)}{2} e^{-2 i \Phi(\tau)} d \tau \tag{5.119}
\end{equation*}
$$

where $\Phi(t)=\int^{t} \sqrt{H_{z}^{2}(t)+H_{x}^{2}(t)} d t$ is the accumulated adiabatic phase, and

$$
\dot{\Theta}(t)=\frac{\left[\dot{H}_{z}(t) H_{z}(t)-\dot{H}_{z}(t) H_{x}(t)\right]}{\left[H_{z}^{2}(t)+H_{x}^{2}(t)\right]}
$$

is the derivative of the angle $\Theta(t)=\pi / 2-\arctan \left(H_{z}(t) / H_{x}(t)\right)$ in the transformation matrix $\exp \left[i \widehat{\sigma}_{y} \Theta(t) / 2\right]$ relating the diabatic and the adiabatic reference frames (5.105). The main contribution to the integral (5.119) comes from the domain around the points $t_{n}$ in the complex plane $t$ given by the relation $\left[H_{z}^{2}\left(t_{n}\right)+H_{x}^{2}\left(t_{n}\right)\right]=0$ where the derivative $\dot{\Theta}(t)$ has poles, and the phase $\Phi(t)$ has a stationary point $\dot{\Phi}(t)=\sqrt{H_{z}^{2}(t)+H_{x}^{2}(t)}=0$. However, the fact that $\dot{\Phi}(t)$ is not analytic at these points does not allow one to employ either the Cauchy formula or the stationary point approximation for $\Phi(t)$. Nevertheless, the Taylor expansion of the square root and the substitution $t_{n}-t=x^{2}$ result in

$$
\begin{equation*}
\psi_{1}^{\prime}(t \rightarrow \infty) \simeq e^{-2 i \Phi\left(t_{n}\right)+i \Phi(t)} \int e^{-i A x} \frac{d x}{2 x}=\pi e^{-2 i \Phi\left(t_{n}\right)+i \Phi(t)} \tag{5.120}
\end{equation*}
$$

and apparently all the contributions of different points $t_{n}$ in that half of the complex plane $t$ where $-2 i \Phi\left(t_{n}\right)$ has a negative real part have to be added. Note that the coefficient $A$ originating from the third-order term in the Taylor expansion does not enter the result. The main contribution corresponds to $t_{n}$ with the largest negative real part.


Fig. 5.8. (a) The integration contour $C$ for the calculation of the non-adiabatic transition probability. (b) Imaginary part of the transition frequency as a function of the complex-valued time. The contour circumvents two complex conjugated branching points $t_{n}=i$ and $t_{n}^{*}=-i$, where the energies corresponding to the state $a$ and the state $b$ coincide, that is $\omega_{a, b}=0$. One has to take either the open contour from $t_{n}^{*}=-i$ to $t_{n}=i$ or take half of the integral over the loop shown.

For the transition probabilities $P_{a b}$ between the adiabatic states $a$ and $b$ the result (5.120) gives a short and elegant generalization known as the Dykhne expression

$$
\begin{equation*}
P_{a b}=\left|\psi_{1}^{\prime}(t \rightarrow \infty)\right| \simeq \pi^{2} \exp \left\{i \int_{t_{n}}^{t_{n}^{*}} \omega_{a b}(t) d t\right\} \tag{5.121}
\end{equation*}
$$

where the integral of the time-dependent energy difference $\omega_{a b}(t)$ between the adiabatic states $a$ and $b$ calculated in the limits suggested by the two closest complex conjugated roots $t_{n}$ and $t_{n}^{*}$ of the equation $\omega_{a b}(t)=0$ takes a positive imaginary value. The other possibility is to integrate half of the frequency over a closed loop, shown in Fig. 5.8.

### 5.3.2 Multilevel Systems

The adiabatic reference system can also be introduced for a multilevel quantum system with a time-dependent Hamiltonian, if all the eigenstates are non-degenerate. Indeed, let $\widehat{H}(t)$ be a time-dependent Hamiltonian entering the Schrödinger equation

$$
\begin{equation*}
i \frac{\partial}{\partial t}|\psi\rangle=\widehat{H}(t)|\psi\rangle \tag{5.122}
\end{equation*}
$$

for a large system of $n+1$ interacting non-degenerate levels corresponding to the states $|k\rangle ; k=0,1, \ldots, n$, The state vector therefore reads $|\psi\rangle=\sum_{k} \psi_{k}|k\rangle$. For each time $t$ one can find a set of eigenvalues $E_{k^{\prime}}(t)$ of the Hamiltonian, and the corresponding eigenvectors $\left|k^{\prime}\right\rangle$, that form an adiabatic basis set related to the initial basis via the expression $|k\rangle=\widehat{O}(t)\left|k^{\prime}\right\rangle$. Here $\widehat{O}(t)$ is a unitary time-dependent operator of the transformation diagonalizing the Hamiltonian, as in (5.105) for the Landau-Zener problem. By analogy, substitution of $|\psi\rangle=\widehat{O}(t)\left|\psi^{\prime}\right\rangle$ into (5.105) yields (5.107) for the state vector $\left|\psi^{\prime}\right\rangle$ in the adiabatic representation and results in the Schrödinger equation (5.109) with the diagonal Hamiltonian $\widehat{H}^{\prime}(t)=$ $\widehat{O}^{-1}(t) \widehat{H}(t) \widehat{O}(t)=\delta_{k^{\prime} m^{\prime}} E_{k^{\prime}}(t)$ perturbed by a small and slowly changing operator $\widehat{C}(t)=\widehat{O}^{-1}(t) \widehat{\dot{O}}(t)$ originating from a slow rotation of the adiabatic basis in Hilbert space. In the adiabatic representation, the last operator does not have diagonal matrix elements.

Note that by passing to the adiabatic reference system, we did not gain much from a formal point of view, since instead of an initial time-dependent Hamiltonian $\widehat{H}(t)$, which does not commute with itself at $t^{\prime} \neq t$, we obtain another one, $\widehat{H}^{\prime}(t)-\widehat{C}(t)$, neither self-commuting for different times. The only advantage is that the off-diagonal elements of the new Hamiltonian are directly related to the rotation of the adiabatic basis, and vanish when the rotation rate approaches zero. When the latter is the case, the transformation $\widehat{O}(t)$ diagonalizing the Hamiltonian also diagonalizes the evolution operator

$$
\begin{equation*}
\widehat{U}_{O}(t)=\widehat{O}(t) \exp \left[-i \int^{t} \widehat{H}^{\prime}(\tau) d \tau\right] \widehat{O}^{-1}(t) \tag{5.123}
\end{equation*}
$$

which implies that no transitions occur among the adiabatic states. On the contrary, in the general case of (5.22) the operator $\widehat{U}(t)$ allows for such nonadiabatic transitions, and therefore is not diagonal in the $\left|k^{\prime}\right\rangle$ representation. This difference originates from the finite $\widehat{C}(t)$ not commuting with $\widehat{H}^{\prime}$ and results in a remainder $\widehat{R}$ encountered in (5.22) which, as we have seen, is exponentially $\left(\sim e^{-\pi V^{2} / \alpha}\right)$ small for the two-level Landau-Zener problem. However, in a multilevel system the large variety of possible paths leading from the initial to the final state may change the situation considerably by augmenting the number of the transition channels possible.

Let us consider this situation in more detail, and calculate the remainder $\widehat{R}$ responsible for the discrepancy between the set of adiabatic states $\left|k^{\prime}\right\rangle$ and the set of eigenvectors of the evolution operator. To get the maximum benefit from the adiabaticity condition, we make use of the interaction representation, that is, we take the evolution operator in the form $\widehat{U}(t)=e^{-i \widehat{R}(t)} \widehat{U}_{O}(t)$. The differential equation for $\widehat{R}(t)$ follows from the equation for the evolution operator $\widehat{U}_{\text {int }}(t)=e^{-i \widehat{R}(t)}$ in the interaction representation

$$
\begin{equation*}
i \frac{\partial}{\partial t} \widehat{U}_{i n t}(t)=\widehat{H}^{\prime \prime}(t) \widehat{U}_{i n t}(t) \tag{5.124}
\end{equation*}
$$

with the Hamiltonian $\widehat{H}^{\prime \prime}(t)=-e^{i \int \widehat{H}^{\prime}(t) d t} \widehat{C}(t) e^{-i \int \widehat{H}^{\prime}(t) d t}$ which in the adiabatic basis contains only small and rapidly oscillating non-diagonal matrix elements. This allows us to find the equation for the coefficients in the from of the Backer-Campbell-Hausdorff formula (5.3). Indeed, from (5.124) one obtains the expression $\widehat{U}(t+d t)=\exp \left[-i \widehat{H}^{\prime \prime}(t) d t\right] \widehat{U}(t)$ for an infinitesimally small variation of the evolution operator, which in agreement with (5.1) yields

$$
\begin{equation*}
e^{-i \widehat{R}(t+d t)}=e^{-i \widehat{H}^{\prime \prime}(t) d t} e^{-i \widehat{R}(t)} \tag{5.125}
\end{equation*}
$$

We take the logarithm of both sides of this equation and expanding the first exponent in $a$ Taylor series obtain

$$
\begin{align*}
-i \widehat{R}(t+d t) & =\ln \left[e^{-i \widehat{R}(t)}-i \widehat{H}^{\prime \prime}(t) e^{-i \widehat{R}(t)} d t\right] \\
& =\ln \left[1-\left(1-e^{-i \widehat{R}(t)}+i \widehat{H}^{\prime \prime}(t) e^{-i \widehat{R}(t)} d t\right)\right] \\
& =\sum_{n} \frac{(-1)^{n+1}}{n}\left[1-e^{-i \widehat{R}(t)}+i \widehat{H}^{\prime \prime}(t) d t e^{-i \widehat{R}(t)}\right]^{n} \tag{5.126}
\end{align*}
$$

Taylor expansion of the logarithm over $\widehat{H}^{\prime \prime}$ yields

$$
\begin{align*}
-i \widehat{R}(t+d t)= & \sum_{n=1} \frac{(-1)^{n+1}}{n}\left[1-e^{-i \widehat{R}(t)}\right]^{n}  \tag{5.127}\\
& +i d t \sum_{n, k=0} \frac{(-1)^{n+k}}{n+k+1}\left[1-e^{-i \widehat{R}(t)}\right]^{n} \widehat{H}^{\prime \prime}(t) e^{-i \widehat{R}(t)}\left[1-e^{-i \widehat{R}(t)}\right]^{k}
\end{align*}
$$

which after performing the summation takes the form

$$
\begin{equation*}
-i \widehat{R}(t+d t)=-i \widehat{R}(t)+i d t \sum_{k=0} a_{k}\left\{[-i \widehat{R}(t),\}^{k} \widehat{H}^{\prime \prime}(t)\{ ]\right\}^{k}, \tag{5.128}
\end{equation*}
$$

where the coefficients $a_{k}$ in the commutator series are given by the Taylor expansion $\sum a_{k} x^{k}=x\left(e^{x}-1\right)^{-1}$. Here we have employed the integral representation $(n+k+1)^{-1}=\int_{0}^{1} z^{n+k} d z$ that enables one to rewrite the sum in the last line of (5.127) in the form

$$
\begin{align*}
& \int_{0}^{1} \sum_{n, k=0}(-z)^{n}\left[1-e^{-i \widehat{R}(t)}\right]^{n} \widehat{H}^{\prime \prime}(t) e^{-i \widehat{R}(t)}(-z)^{k}\left[1-e^{-i \widehat{R}(t)}\right]^{k} d z \\
= & \int_{0}^{1}\left[1+z\left(1-e^{-i \widehat{R}(t)}\right)\right]^{-1} \widehat{H}^{\prime \prime}(t) e^{-i \widehat{R}(t)}\left[1+z\left(1-e^{-i \widehat{R}(t)}\right)\right]^{-1} d z . \tag{5.129}
\end{align*}
$$

For the matrix elements of (5.129) in the representation where $\widehat{R}(t)$ is diagonal we have performed the integration

$$
\begin{align*}
& \int_{0}^{1} \frac{1}{1+z\left(1-e^{-i \widehat{R}_{l l}(t)}\right)} \widehat{H}_{l m}^{\prime \prime}(t) \frac{e^{-i \widehat{R}_{m m}(t)}}{1+z\left(1-e^{-i \widehat{R}_{m m}(t)}\right)} d z \\
= & \frac{\log \left[e^{-i \widehat{R}_{l l}(t)}\right]-\log \left[e^{-i \widehat{R}_{m m}(t)}\right]}{e^{-i \widehat{R}_{l l}(t)}-e^{-i \widehat{R}_{m m}(t)}} e^{-i \widehat{R}_{m m}(t)} \widehat{H}_{l m}^{\prime \prime}(t) \\
= & \sum_{k=0} a_{k}\left[i \widehat{R}_{m m}(t)-i \widehat{R}_{l l}(t)\right]^{k} \widehat{H}_{l m}^{\prime \prime}(t), \tag{5.130}
\end{align*}
$$

which yields the matrix form of the commutator expansion (5.128).
The nonlinear differential equation (5.128)

$$
\begin{equation*}
\frac{\partial \widehat{R}(t)}{\partial t}=-\sum_{k=0} a_{k}\left\{[-i \widehat{R}(t),\}^{k} \widehat{H}^{\prime \prime}(t)\{ ]\right\}^{k} \tag{5.131}
\end{equation*}
$$

for the matrix $\widehat{R}(t)$ has, in the adiabatic case, an approximate solution based on the fact that $\widehat{H}^{\prime \prime}(t)$ is small and rapidly changing. Note that assuming $\widehat{H}^{\prime \prime}\left(t>t_{A}\right)=\widehat{H}_{B}$ and $\widehat{R}\left(t_{A}\right)=t_{A} \widehat{H}_{A}$ and solving this equation iteratively one can also determine the coefficients $c_{k, s}$ in the Baker-Campbell-Hausdorff formula (5.3)-(5.4).

For small and rapidly changing $\widehat{H}^{\prime \prime}(t)$, by taking only the first term on the right-hand side of (5.131) one obtains $\widehat{R}(t)=\int_{-\infty}^{t} \widehat{H}^{\prime \prime}(t) d t+\widehat{o}(t)$ where the small higher-order correction $\widehat{o}(t)$ can also be found by iterations if necessary. Neglecting the correction immediately results in

$$
\begin{equation*}
\widehat{R}(\infty)=-\int_{-\infty}^{\infty} e^{i \int^{t} \widehat{H}^{\prime}(\tau) d \tau} \widehat{O}^{-1}(t) \widehat{\dot{O}}(t) e^{-i \int^{t} \widehat{H}^{\prime}(\tau) d \tau} d t \tag{5.132}
\end{equation*}
$$

where we have substituted the explicit form of $\widehat{H}^{\prime \prime}(t)$. Matrix elements

$$
\begin{equation*}
R_{l m}(\infty)=-\sum_{k} \int_{-\infty}^{\infty} e^{i \int^{t} \widehat{H}_{l l}^{\prime}(\tau) d \tau}\left[\widehat{O}^{-1}(t)\right]_{l k} \dot{O}_{k m}(t) e^{-i \int^{t} \widehat{H}_{m m}^{\prime}(\tau) d \tau} d t \tag{5.133}
\end{equation*}
$$

found by an approximate calculation similar to (5.121) give the amplitudes of the direct transitions between the states $l$ and $m$ in the presence of a number of other states $k$, whereas the overall probabilities that include sequential transitions in this order of approximation are given by the expression

$$
\begin{equation*}
P_{l m}=\left|\left[e^{-i \widehat{R}(\infty)}\right]_{l m}\right|^{2} \tag{5.134}
\end{equation*}
$$

In the adiabatic regime each matrix element $R_{k n}$ of $\widehat{R}(\infty)$ is exponentially small. However, this does not necessarily mean that the matrix itself is small, since the rank $\sim N$ of the matrix and hence its norm $\sim \sqrt{N\left(R_{k n}\right)^{2}}$ can be large for a large number $N$ of levels in the system. Therefore, the Taylor expansion of the exponent in (5.134) cannot be automatically truncated at the first-order terms. The question arises whether equation (5.131) for $\widehat{R}(t)$ can be restricted to the zero-order approximation in such a situation that has given the result (5.132)? The answer can be positive if the system is "complex enough", such that both the Hamiltonian $\widehat{H}^{\prime \prime}(t)$ and its primitive $\int \widehat{H}^{\prime \prime}(t) d t$ are generic, that is each of them contains all or almost all of the $\sim N^{2}$ terms in their representations in the form of linear combinations of the operators $\widehat{\mathcal{G}}_{n}$ generating the Lie group of the $N \times N$ matrix of evolution. Among the total number $N^{2} \times N^{2}$ of pairs of generators the number of noncommuting pairs is of the order of $N^{2}$, and hence a typical matrix element of the commutator of two generic matrices is of the order of the product of their typical matrix elements. Therefore the norms of commutators of small generic matrices are small values of a higher order, and thus the correction $\widehat{o}(t)=\widehat{R}(t)-\int_{-\infty}^{t} \widehat{H}^{\prime \prime}(t) d t$ can indeed be neglected.

### 5.4 Time-Dependent Level-Band System

Earlierin Chap. 3, we have seen that a single level interacting with a homogeneous band of levels is a very particular example of complex multilevel quantum systems that allow an exact analytic description. A number of problems for the time-dependent level-band systems with the Schrödinger equation

$$
\begin{align*}
i \dot{\psi}_{n} & =\delta n \psi_{n}+V_{n 0}(t) \psi_{0} \\
i \dot{\psi}_{0} & =\Delta(t) \psi_{0}+\sum_{n=1}^{\infty} V_{0 n}(t) \psi_{n} \tag{5.135}
\end{align*}
$$

can also be solved analytically for certain particular cases of the time dependence of the detuning $\Delta(t)$ and the couplings $V_{n 0}(t)$. Here we consider just some of them, which manifest meaningful and qualitatively important features.

One peculiarity, common to all of the problems under consideration, has to be mentioned in advance: we assume that the size of the band $X$ is very large, much larger, as compared to any of the energy parameters of the system. However we cannot simply set it to infinity since the logarithm of the ratio $X / \delta$, that is the logarithm of the typical number of levels in the system, remains one of the crucial parameters of the problem. It has a very moderate size even when the bandwidth $X$ is enormous, and enters many of the characteristics of the system, such as the typical response time.

### 5.4.1 The Demkov-Osherov Problem

Let us consider now an example of complex multilevel quantum systems with time-dependent parameters that have an exact analytic solution. A system shown in Fig. 5.9 composed of a non-degenerate quantum state $|0\rangle$ with a time-dependent energy $E_{0}(t)=\alpha t$ interacting with an equidistant spectrum $E_{n}=\delta n$ via a constant coupling $V_{0 n}=V$ was suggested by Demkov and Osherov as a model for chemical reactions. The Schrödinger equation (5.135) for such a system

$$
\begin{align*}
i \dot{\psi}_{n} & =\delta n \psi_{n}+V \psi_{0} \\
i \dot{\psi}_{0} & =\alpha t \psi_{0}+V \sum_{n=1}^{\infty} \psi_{n} \tag{5.136}
\end{align*}
$$

is similar to (3.23) for the level-band system of Sect. 3.2 with the only difference that the energy of the initially populated state $|0\rangle$ depends linearly on time. The Laplace contour integral method (see p. 205) allows one to find an exact solution. We look for the solution in the form

$$
\begin{equation*}
\psi_{n}(t)=\int_{C} \psi_{n}(\varepsilon) e^{-i \varepsilon t} d \varepsilon, \tag{5.137}
\end{equation*}
$$

where the integration contour $C$ is selected in such a way that $\psi_{n}(\varepsilon)$ assumes the same values at the ends of $C$ for any $n=0,1, \ldots$ while the amplitudes themselves satisfy the initial conditions. This choice allows one to perform integration in parts, replacing the time derivatives $\partial / \partial t$ by the factors $-i \varepsilon$ and factors $t$ by $-i \partial / \partial \varepsilon$, which results in a set of equations for the amplitudes entering the integrands

$$
\begin{align*}
& \varepsilon \psi_{n}(\varepsilon)=\delta n \psi_{n}(\varepsilon)+V \psi_{0}(\varepsilon) \\
& \varepsilon \psi_{0}(\varepsilon)=-i \alpha \frac{\partial \psi_{0}(\varepsilon)}{\partial \varepsilon}+V \sum_{n=1}^{\infty} \psi_{n}(\varepsilon) . \tag{5.138}
\end{align*}
$$



Fig. 5.9. (a) A level with energy linearly dependent on time crosses a band of levels. The solid and the dashed lines represent the adiabatic and the diabatic state energies, respectively. (b) Integration contours for different initial conditions. For the moving state initially populated one has to take the contour $C$, while for the initial condition $\psi_{n}(t=-\infty)=1$ the contour $C^{\prime}$ has to be selected. The borders of the Stokes domains are shown by the dashed lines.

We note here an important difference between (5.138) and equations resulting from the standard Fourier-Laplace method of solving systems of linear differential equations discussed earlier in Chap. 3: this equation does not explicitly contain initial conditions, and therefore one has to look for a proper integration contour instead.

For $\psi_{0}(\varepsilon)$ (5.138) immediately yields a first-order differential equation

$$
\begin{equation*}
i \alpha \frac{\partial \psi_{0}}{\partial \varepsilon}+\varepsilon \psi_{0}-\sum_{n=1}^{\infty} \frac{V^{2} \psi_{0}}{\varepsilon-\delta n}=0 \tag{5.139}
\end{equation*}
$$

with the solution

$$
\begin{align*}
\psi_{0} & =\exp \left\{\int\left(i \frac{\varepsilon}{\alpha}-\sum_{n=1}^{\infty} \frac{i V^{2} / \alpha}{\varepsilon-\delta n}\right) d \varepsilon\right\} \\
& =\exp \left\{i \frac{\varepsilon^{2}}{2 \alpha}-\sum_{n=1}^{\infty} \frac{i V^{2}}{\alpha} \ln \left(\frac{\varepsilon-\delta n}{A}\right)\right\}, \tag{5.140}
\end{align*}
$$

where $A$ is an integration constant. Substitution of this expression into (5.137) yields

$$
\begin{equation*}
\psi_{0}(t)=\int_{C} \exp \left\{i \frac{\varepsilon^{2}}{2 \alpha}-i \varepsilon t-\sum_{n=1}^{\infty} \frac{i V^{2}}{\alpha} \ln \left(\frac{\varepsilon-\delta n}{A}\right)\right\} d \varepsilon \tag{5.141}
\end{equation*}
$$

Now we have to find an integration contour that would correspond to the initial conditions imposed. The expression in the exponent has no poles,
but only branching points at each $\varepsilon=\delta n$, such that closed loops are incapable of giving non-trivial solutions, while the possibility to find an open loop with a pair of initial and final points where the integrand (5.137) takes identical values for each $n$ is quite questionable. The other possibility is to take contours where all of the integrands vanish. This can be done either for $\varepsilon \rightarrow \pm i e^{ \pm i \pi / 4} \infty$ or for $\arg (\varepsilon-\delta n) \rightarrow-\infty$. These two possibilities shown in Fig. 5.9 correspond to different initial conditions. If all of the population was initially in the moving level, the contour $C$ going from $\varepsilon=\sqrt{i} \infty$ to $\varepsilon=-\sqrt{i} \infty$ along the Stokes lines can be selected. At $t \sim t_{0} \rightarrow \infty$ where the functional dependence of the logarithmic terms is negligible, the contour passes through the saddle point $\varepsilon=\alpha t$ on the real axis, which yields

$$
\begin{equation*}
\psi_{0}(t)=\text { const } \int_{C} \exp \left\{i \frac{\varepsilon^{2}}{2 \alpha}-i \varepsilon t\right\} d \varepsilon=\exp \left\{-i \frac{t^{2} \alpha}{2}\right\} \tag{5.142}
\end{equation*}
$$

The contour $C^{\prime}$ appropriate for the initial condition $\left.\psi_{n}(t)\right|_{t \rightarrow-\infty}=1$ is a spiral winding towards the point $\varepsilon=\delta n$, and the main contribution for large negative $t$ comes from the vicinity of the saddle point $\varepsilon=\delta n-\left(V^{2} / \alpha+i\right) / t$

$$
\begin{align*}
\psi_{n}(t) & =\int_{C^{\prime}} \frac{V}{\varepsilon-\delta n} \exp \left\{i \frac{\varepsilon^{2}}{2 \alpha}-i \varepsilon t-\sum_{n=1}^{\infty} \frac{i V^{2}}{\alpha} \ln \left(\frac{\varepsilon-\delta n}{A}\right)\right\} d \varepsilon \\
& \simeq \exp \left\{-i \delta n t+\frac{i V^{2}}{\alpha} \ln \left(-t / A^{\prime}\right)\right\} . \tag{5.143}
\end{align*}
$$

where $A^{\prime}$ is a constant different from $A$.
Let us consider the problem in the limit $\delta \rightarrow 0, V^{2} / \delta \rightarrow W$, and for $t \ll 1 / \delta$, namely assuming that our time interval $t$ being long, still remains shorter than the Heisenberg recurrence time given by the inverse of the level spacing $\delta$. Replacing the sum in (5.140) by the integral we arrive at

$$
\begin{equation*}
\psi_{0}=\exp \left\{i \frac{\varepsilon^{2}}{2 \alpha}-\frac{i W}{\alpha} \int \ln \left(\frac{\varepsilon-x}{A}\right) d x\right\} \tag{5.144}
\end{equation*}
$$

Evaluation of the integral in the finite limits $0<x<X<\infty$ results in

$$
\begin{align*}
\int_{0}^{X} \ln \left(\frac{\varepsilon-x}{A}\right) d x & =\varepsilon \ln \left(\frac{-\varepsilon}{X-\varepsilon}\right)-X+X \ln \left[\frac{-X+\varepsilon}{A}\right] \\
& \simeq \varepsilon \ln \left(\frac{-\varepsilon}{X}\right)-\varepsilon+\text { const }=\varepsilon \log \left(\frac{-\varepsilon}{X e}\right)+\text { const } \tag{5.145}
\end{align*}
$$

where we have ignored the small variable $\varepsilon$ in the arguments of logarithms as compared to a large constant $X$ in the first term on the right-hand side and retain only the correction linear in $\varepsilon$ resulting from the third term; these are the only two non-constant parts in this extreme. Equation (5.144) then takes the form

$$
\begin{equation*}
\psi_{0}=\text { const } \int \exp \left\{i \frac{\varepsilon^{2}}{2 \alpha}-\frac{i W}{\alpha} \varepsilon \ln \left(\frac{-\varepsilon}{X e}\right)\right\} d \varepsilon \tag{5.146}
\end{equation*}
$$

and for the time-dependent amplitudes results in

$$
\begin{equation*}
\psi_{0}=\frac{1}{\sqrt{2 \alpha \pi}} \int_{C} \exp \left\{i \frac{\varepsilon^{2}}{2 \alpha}+i\left(\frac{W}{\alpha} \ln \frac{X}{W}-t\right) \varepsilon-\varepsilon \frac{i W}{\alpha} \ln \left(\frac{-\varepsilon}{W e}\right)\right\} d \varepsilon \tag{5.147}
\end{equation*}
$$

where we have employed the relation $\log (-\varepsilon / X)=\ln (-\varepsilon / W)-\ln (X / W)$, and have chosen the constant corresponding to the initial condition $\left|\psi_{0}\right|_{t \rightarrow-\infty} \mid=$ 1. By the replacement $\varepsilon \rightarrow y W$ and $t=(W / \alpha) \ln (X / W)+W T / \alpha$ we express the probability amplitude in terms of a contour integral

$$
\begin{equation*}
\psi_{0}=\sqrt{\frac{W^{2}}{2 \alpha \pi}} \int_{C} \exp \left\{\frac{i W^{2}}{\alpha}\left[\frac{y^{2}}{2}-T y-y \ln \left(\frac{-y}{e}\right)\right]\right\} d y \tag{5.148}
\end{equation*}
$$

which is a universal function of two parameters, $T$ and $W^{2} / \alpha$. In Fig. 5.10 we depict $\left|\psi_{0}\right|$ as a function of these parameters. One sees that with the elapse of time the population $\left|\psi_{0}\right|^{2}$ initially equal to 1 approaches zero with a rate proportional to the parameter $W^{2} / \alpha$. This process is not just an exponential relaxation - a tiny reminiscence of the coherence manifests itself in small oscillations of the real and imaginary parts of the amplitude. The asymptotic expression for the probability of the state $|0\rangle$ reads

$$
\begin{equation*}
\rho_{0} \simeq e^{-2 \pi W^{2} T / \alpha} \tag{5.149}
\end{equation*}
$$

Equation (5.149) written in the ordinary time variable $t \sim T W / \alpha$ corresponds to the exponential decay $\rho_{0} \sim e^{-2 \pi V^{2} t / \delta}$ identical to that in the case of a stationary level-band system, which is given by (3.34).

In Sect. 5.1 we have noticed that the complexity of the multilevel quantum system dynamics originates from the algebraic properties of the Hamiltonian and its time derivatives. If the operator algebra of their commutators is rich enough, the dynamics of the system does not have a simple description either in terms of explicit algebraic expressions nor in terms of integrals. The fact that the Demkov-Osherov problem under consideration possesses such an explicit solution implies that it corresponds to a very particular structure of the underlying operator relations. Inspection immediately shows that only the first derivative of the Hamiltonian corresponding to (5.136) is a nonzero matrix of rank 1 , and is proportional to the operator $\widehat{P}_{0}$ of projection to the moving state $|0\rangle$. Nevertheless, the direct application of the idea of algebraic description is not of much use, since the ensemble of all possible commutations of the Hamiltonian $\left.i \widehat{H}(t)\right|_{t=0}$ and the projector $i \widehat{P}_{0}$ spans all operator space. In other words, any Hermitian operator in the Hilbert space of quantum states can be given in the form of a linear combination of $i \widehat{H}(0), i \widehat{P}_{0}$, and their commutators of first and all higher orders, which


Fig. 5.10. Probability amplitude to remain in the initial state $|\psi|$ as a function of time $T=\alpha t / W-\ln (X / W)$ and the interaction parameter $W^{2} / \alpha$ for the linearly moving level coupled to a continuum. (a) The amplitudes tends from its initial value 1 to the final value 0 with the rate depending on the scaled coupling $\sqrt{W^{2} / \alpha}$. (b) Manifestation of a coherent transitional process that takes place at the moment that the level enters the continuum. (c) Real (dashed line) and imaginary (dot-dashed line) parts of the amplitude during the transitional process, for the value of $W^{2} / \alpha$ shown in (b) by the arrow.
form therefore a complete basis set. Hence from the viewpoint of the algebra of operators in the Hilbert space of quantum states, we have encountered a problem of the general type.

The possibility of finding the analytical solution comes from the fact that the relation $\partial \widehat{H}(t) /\left.\partial t\right|_{t=0} \sim \widehat{P}_{0}$ results in a very important simplification of the problem. After the Fourier transformation $t \rightarrow \varepsilon$ the Schrödinger equation results in a linear first-order differential equation (5.139) in the space $\varepsilon$. In other words, the algebraic structure of the problem is very simple from the viewpoint of operators in the frequency space. Indeed, the commutator of the operator $\partial / \partial \varepsilon$ canonically conjugated with $\varepsilon$ and any multiplicative operator $U(\varepsilon)$ results in another multiplicative operator $U^{\prime}(\varepsilon)=[\partial / \partial \varepsilon, U(\varepsilon)]$ which is equivalent to $U(\varepsilon)$ as an operator. We note that any further complication of the problem, such as parabolic time dependence of the detuning, or simultaneous linear time dependence of two states, results, instead of (5.139), in a second-order differential equation for $\psi_{0}$, whose algebraic structure is much more involved, and which in the general case does not have an explicit analytic solution.

### 5.4.2 The Landau-Zener Transition at the Continuum Edge

When considering on p. 74 the role of spectrum edges, we noticed that a discrete quantum state exists near the end of the continuum interacting with a non-degenerate level $|0\rangle$. The positions of such states in the vicinity of each of two continuum edges are found from the solution of a transcendental equation illustrated in Fig. 3.6, and their population remains constant with the course of time. Therefore, even though weakly populated, these states bring the dominating contribution to the long-time asymptotic of the level population, while the contribution of the rest of the spectrum decays exponentially. For the level-band problem with time-dependent parameters such discrete eigenvalues of instantaneous Hamiltonians exist as well and although they move, their population also plays an important role in the long time dynamics of the system.

## The single moving level

One can trace the origin of these discrete states for the Demkov-Osherov problem in the limit $\delta \rightarrow 0, V^{2} / \delta \rightarrow W$ shown in Fig. 5.11(a). We concentrate


Fig. 5.11. Level crosses a continuous band (a). There are two discrete adiabatic states above and below the continuum. Population (b) of the lower adiabatic state as a function of the scaled time $T=\alpha t / V^{2} g-\ln X / W$ and the adiabaticity parameter $W / \sqrt{\alpha}$.
here on the role of the lower edge, considering the problem only as long as the level remains far below the upper edge of the band. Even in the condition where the level $|0\rangle$ crosses the lower edge and penetrates deeper and deeper into the band, the adiabatic state $|a d\rangle$ corresponding at $t \rightarrow-\infty$ to the position $E_{a d}(t) \rightarrow \alpha t$ of the level $|0\rangle$ remains discrete, although it approaches the lower edge of the continuum, that is $E_{a d}(t) \rightarrow 0$ at $t \rightarrow \infty$. Our aim is to find the time-dependent population of this adiabatic state.

Let us first express the probability amplitude of the adiabatic state in terms of the amplitudes $\psi_{n}$ of diabatic states $|n\rangle$ of (5.136) and than take the limit $\delta \rightarrow 0$. From the equation for the eigenvalues one finds immediately

$$
\begin{align*}
& E_{a d}(t) \widetilde{\psi}_{n}=\delta n \widetilde{\psi}_{n}+V \widetilde{\psi}_{0} \\
& E_{a d}(t) \widetilde{\psi}_{0}=\alpha t \widetilde{\psi}_{0}+V \sum_{n=1}^{\infty} \widetilde{\psi}_{n} \tag{5.150}
\end{align*}
$$

which yields algebraic equations for the projections $\langle a d \mid n\rangle=\widetilde{\psi}_{n}$ and for the adiabatic energy

$$
\begin{align*}
\widetilde{\psi}_{n} & =\frac{V}{E_{a d}(t)-\delta n} \widetilde{\psi}_{0} \\
E_{a d}(t) & =\alpha t+\sum_{n=1}^{\infty} \frac{V^{2}}{E_{a d}(t)-\delta n}, \tag{5.151}
\end{align*}
$$

respectively. The first equation allows one to express the normalized adiabatic state in the form

$$
\begin{equation*}
|a d\rangle=\frac{1}{\mathcal{N}}\left(|0\rangle+\sum_{n=1}^{\infty} \frac{V}{E_{a d}(t)-\delta n}|n\rangle\right) \tag{5.152}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{N}=\sqrt{1+\sum_{n=1}^{\infty}\left(\frac{V}{E_{a d}(t)-\delta n}\right)^{2}} \tag{5.153}
\end{equation*}
$$

and therefore relates the amplitude $\psi_{a d}=\langle a d \mid s t\rangle$ of the adiabatic state to the solutions $\psi_{n}(t)$ of (5.136) via the standard definition of the state vector $|s t\rangle=\sum \psi_{n}(t)|n\rangle$

$$
\begin{align*}
\psi_{a d}(t) & =\frac{1}{\mathcal{N}}\left(\langle 0|+\sum_{n=1}^{\infty} \frac{V}{E_{\text {ad }}(t)-\delta n}\langle n|\right) \sum_{n=0}^{\infty} \psi_{n}(t)|n\rangle \\
& =\frac{1}{\mathcal{N}}\left(\psi_{0}+\sum_{n=1}^{\infty} \frac{V \psi_{n}}{E_{\text {ad }}(t)-\delta n}\right) . \tag{5.154}
\end{align*}
$$

The first equation in the set (5.136) yields

$$
\begin{equation*}
\psi_{n}(t)=-i V \int_{-\infty}^{t} e^{-i \delta n\left(t-t^{\prime}\right)} \psi_{0}\left(t^{\prime}\right) d t \tag{5.155}
\end{equation*}
$$

and hence

$$
\begin{equation*}
\psi_{a d}(t)=\frac{1}{\mathcal{N}}\left(\psi_{0}(t)-i \int_{-\infty}^{t} \sum_{n=1}^{\infty} \frac{V^{2} e^{-i \delta n\left(t-t^{\prime}\right)} \psi_{0}\left(t^{\prime}\right)}{E_{a d}(t)-\delta n} d t\right) \tag{5.156}
\end{equation*}
$$

where the probability amplitude $\psi_{0}(t)$ to remain in the state $|0\rangle$ is given by (5.141).

Now we are in a position to consider the limit $\delta \rightarrow 0$ and replace summation by integration. This yields an explicit expression for the probability amplitude

$$
\begin{equation*}
\psi_{a d}(t)=\frac{\psi_{0}(t)-i \int_{-\infty}^{t} W \int_{0}^{\infty} \frac{\exp \left[-i x\left(t-t^{\prime}\right)\right]}{E_{a d}(t)-x} d x \psi_{0}\left(t^{\prime}\right) d t^{\prime}}{\sqrt{1-W / E_{a d}(t)}} \tag{5.157}
\end{equation*}
$$

where we have taken into account the normalization constant

$$
\begin{equation*}
\mathcal{N}=\sqrt{1+\int_{0}^{\infty} \frac{W d x}{\left(E_{a d}(t)-x\right)^{2}}}=\sqrt{1-\frac{W}{E_{a d}(t)}} \tag{5.158}
\end{equation*}
$$

At the same extreme equation (5.151) reads

$$
\begin{equation*}
E_{a d}(t)=\alpha t+\int_{0}^{X \rightarrow \infty} \frac{W d x}{E_{a d}(t)-x}=\alpha t+W \ln \left(\frac{E_{a d}(t)}{-X}\right) \tag{5.159}
\end{equation*}
$$

where we should take precautions and cut the logarithmic divergence at the upper integration limit, by choosing a large but finite width $X$ for the continuum band in the same way as earlier in (5.145)-(5.148). As then, the replacement $t \rightarrow(W / \alpha) \log (X / W)+W T / \alpha$, along with the substitutions $E_{a d}(t)=\mathcal{E}(T) W, W^{2} / \alpha=w$ allows one to get rid of the large constant $X$, and after a proper change of the integration variable $x$ together with (5.148) yields

$$
\begin{align*}
\psi_{a d} & =\frac{1}{\sqrt{1-\mathcal{E}^{-1}}}\left(\psi_{0}-i \int_{-\infty}^{T} w \int_{0}^{\infty} \frac{\exp \left[-i x\left(T-T^{\prime}\right)\right]}{w \mathcal{E}-x} \psi_{0}\left(T^{\prime}\right) d x d T^{\prime}\right) \\
\mathcal{E} & =T+\ln [-\mathcal{E}] \\
\psi_{0} & =\sqrt{\frac{w}{2 \pi}} \int_{C} \exp \left\{i w\left[\frac{y^{2}}{2}-T y-y \ln \left(\frac{-y}{e}\right)\right]\right\} d y \tag{5.160}
\end{align*}
$$

where we have to keep in mind that $\mathcal{E}=\mathcal{E}(T)<0$, where the integration contour $C$ is shown in Fig. 5.9(b).

Substitution of the last equation of set (5.160) into the second term in the parentheses of the first equation followed by integration over $T^{\prime}$ yields

$$
\begin{equation*}
-\sqrt{\frac{w^{3}}{2 \pi}} \int_{C} \int_{0}^{\infty} \frac{\exp \left\{i w\left[\frac{y^{2}}{2}-T y-y \ln \left(\frac{-y}{e}\right)\right]\right\}}{(w \mathcal{E}(T)-x)(i x-i w y)} d x d y \tag{5.161}
\end{equation*}
$$

We note that the position $\mathcal{E}(T)$ of the discrete adiabatic level underneath the band is related to time by the second equation (5.160), and therefore the replacement $T \rightarrow \mathcal{E}-\ln [-\mathcal{E}]$ can be performed in order to simplify the integrals. After integration of (5.161) over $x$ and substitution of the result into the first equation (5.160) we arrive at

$$
\begin{equation*}
\psi_{a d}=\sqrt{\frac{w}{2 \pi}} \int_{C}\left(1+\frac{\ln \frac{-\mathcal{E}}{y}}{y-\mathcal{E}}\right) \frac{\exp \left\{i w\left[\frac{y^{2}}{2}-\mathcal{E} y-y \ln \left(\frac{y}{e \mathcal{E}}\right)\right]\right\}}{\sqrt{\mathcal{E}^{-1}-1}} d y \tag{5.162}
\end{equation*}
$$

In the exponent one recognizes the primitive of $y-\mathcal{E}+\log (-\mathcal{E} / y)$ which is also the numerator of the prefactor, and therefore after integration by parts, the integral takes the form

$$
\begin{equation*}
\psi_{a d}=\sqrt{\frac{w}{2 \pi}} \int_{C} \frac{\exp \left\{i w\left[\frac{y^{2}}{2}-\mathcal{E} y-y \ln \left(\frac{y}{e \mathcal{E}}\right)\right]\right\}}{(y-\mathcal{E})^{2} \sqrt{\mathcal{E}^{-1}-1}} d y \tag{5.163}
\end{equation*}
$$

convenient for numerical calculations.
In Fig. 5.11(b) we show the results of the numerical calculation of the probability $\left|\psi_{a d}\right|^{2}$ given by (5.163) as a function of two parameters $T$ and $w$ where for large couplings and slow motion, one sees that the population decay of the adiabatic state at $T \rightarrow \infty$ is slower than for the state $\psi_{0}$. We demonstrate this in Fig. 5.12, where the adiabatic state population, $\rho_{a d}=$ $\left|\psi_{a d}\right|^{2}$, and its ratio to the diabatic state population $\rho_{0}=\left|\psi_{0}\right|^{2}$ are shown as a function of the parameter $W / \sqrt{\alpha}$ and the energy position $\mathcal{E}(T)$ of the adiabatic level which serves as a nonlinear time-scale.


Fig. 5.12. Population (a) of the adiabatic state as a function of the adiabaticity parameter $W / \alpha^{1 / 2}$ and the position $\mathcal{E}(T)$ of the adiabatic level, given by (5.160) and its ratio (b) to the diabatic state population as a function of the same variables.

As we have mentioned in the previous section, the amplitude $\psi_{0}$ exponentially decreases in the long-time limit such that the population given by
(5.149) decreases as $e^{-2 \pi W t}$, and in a certain regime its asymptotic contribution to the amplitude of the adiabatic state can become small as compared to the contribution of the ensemble of the band states, the states near the edge in particular. In order to find the asymptotic behavior of the adiabatic population, we evaluate the integral (5.163) approximately for large $w$ and long times (that is small $\mathcal{E}$ ). We transform the integration contour to a loop circumventing the branching point 0 and consider the sum of two parts of the integral as the single integral

$$
\begin{equation*}
\psi_{a d}=\int_{-\infty}^{0}\left(1-e^{2 \pi w y}\right) \frac{\exp \left\{i w\left[\frac{y^{2}}{2}-\mathcal{E} y-y \ln \left(\frac{y}{e \mathcal{E}}\right)\right]\right\}}{(y-\mathcal{E})^{2} \sqrt{2 \pi / w} \sqrt{\mathcal{E}^{-1}-1}} d y \tag{5.164}
\end{equation*}
$$

We now set $\left(1-e^{2 \pi w y}\right) \simeq 2 \pi w y$ and keep only the leading term in the exponent, which yields

$$
\begin{align*}
\psi_{a d} & =\frac{-\sqrt{2 \pi w^{3}}}{\sqrt{\mathcal{E}^{-1}-1}} \int_{-\infty}^{0} y \frac{\exp \{i w y \ln (-\mathcal{E})\}}{(y-\mathcal{E})^{2}} d y \\
& \simeq \frac{-\ln (-w \mathcal{E} \ln (-\mathcal{E}))}{\sqrt{\mathcal{E}^{-1}-1}} \sqrt{2 \pi w^{3}} \tag{5.165}
\end{align*}
$$

We note that at long times $T \simeq-\ln [-\mathcal{E}]$ and hence $\mathcal{E} \simeq-e^{-T}$, and therefore

$$
\begin{align*}
\left|\psi_{a d}\right|^{2} & \simeq 2 \pi w^{3} e^{-T}\left[\ln \left(w e^{-T} \ln \left(e^{-T}\right)\right)\right]^{2} \\
& \simeq 2 \pi w^{3} e^{-T} T^{2} \tag{5.166}
\end{align*}
$$

It is expedient to compare the population (5.166) of the adiabatic state with the population of the diabatic state $|0\rangle$ of (5.149). The comparison

$$
\begin{align*}
\rho_{a d} & \simeq 2 \pi w^{3} e^{-T} T^{2} \\
\rho_{0} & \simeq e^{-2 \pi w T}, \tag{5.167}
\end{align*}
$$

shows clearly that for a slow motion $(w \rightarrow \infty)$, the adiabatic contribution dominates. Both expressions (5.167) contain implicitly the width of the band, which determines, however, just the typical time $T$ when the asymptotic regime is attained. Comparison of (5.149) with (5.166) shows that the adiabatic asymptotic regime certainly becomes more important at $T>\log \left(2 \pi w^{3}\right) /(1-2 \pi w)$.

Note another important consequence. Imagine that the level $|0\rangle$ enters the band and remains for a while before returning back, as shown in Fig. 5.13(a). Let us assume that during the return, the system has lost the phase coherence such that the probability to return back to the initial state after the double passage of the continuum edge is just the product of two equal probabilities for the forward and for the back transition, no matter whether we consider the adiabatic or the diabatic trajectory. As for the overall probability $\rho_{0 \rightarrow 0}$



Fig. 5.13. Double passage of the continuum edge (a). The diabatic (dashed line) and the adiabatic (solid line) channels. Population (b) of the non-adiabatic transition as a function of the level velocity $v=\alpha / W^{2}$ for $W=10^{-2}$. The interaction time is taken such that $\mathcal{E}\left(T_{\max }\right) / W=1 / e$
to return in the initial state $|0\rangle$ at $t \rightarrow \infty$ one has to take into account the fact that the adiabatic and the diabatic discrete states are not orthogonal, and therefore this probability has to be taken in the form

$$
\begin{equation*}
\rho_{0 \rightarrow 0} \simeq \rho_{0}^{2}+\rho_{a d}^{2}-\frac{\rho_{0} \rho_{a d}}{|\langle 0 \mid a d\rangle|^{2}}=\rho_{0}^{2}+\rho_{a d}^{2}-\frac{\rho_{0} \rho_{a d}}{|\mathcal{N}|^{2}} \tag{5.168}
\end{equation*}
$$

where the normalization factor $\mathcal{N}$ is given by (5.153). The irreversible transition probability $\rho_{t r}=1-\rho_{0 \rightarrow 0}$ as a function of the velocity has a maximum when the adiabatic and the diabatic populations are of the same order of magnitude, which is illustrated in Fig.5.13(b).

## The Exponentially Rising Interaction

Earlier, on p. 210 we considered the two-level system with an exponentially rising coupling. Here we generalize this consideration to the case of the levelband system. The Schrödinger equation (5.135) for this case reads

$$
\begin{align*}
i \dot{\psi}_{n} & =n \delta \psi_{n}+V e^{\alpha t} \psi_{0} \\
i \dot{\psi}_{0} & =\Delta_{0} \psi_{0}+V e^{\alpha t} \sum_{n=1}^{\infty} \psi_{n} \tag{5.169}
\end{align*}
$$

and with the notation $\theta=\sqrt{\alpha / 2 \delta} \int^{t} V e^{\alpha x} d x=V e^{\alpha t} / \sqrt{2 \alpha \delta}$ and $\Delta=\Delta_{0} / 2 \alpha$, similar to that introduced earlier in (5.94) and (5.85), it adopts the form

$$
\begin{align*}
i \frac{\partial \psi_{n}}{\partial \theta} & =\frac{\Delta+n \delta^{\prime}}{\theta} \psi_{n}+\sqrt{2 \delta^{\prime}} \psi_{0} \\
i \frac{\partial \psi_{0}}{\partial \theta} & =\sqrt{2 \delta^{\prime}} \sum_{n=1}^{\infty} \psi_{n} \tag{5.170}
\end{align*}
$$

resembling (5.95) for a two-level system. Here $\delta^{\prime}=\delta / \alpha$ and the initial condition $\psi_{0}(\theta=0)=1$ is the same as earlier.

Now one can employ the same method and look for the solution in the form of a power series

$$
\begin{align*}
& \psi_{0}(\theta)=\sum_{k=0}^{\infty}(-i)^{k} \frac{c_{0, k}}{k!} \theta^{2 k} \\
& \psi_{n}(\theta)=\sum_{k=0}^{\infty}(-i)^{k} \frac{c_{n, k}}{k!} \theta^{2 k+1} \tag{5.171}
\end{align*}
$$

as was done for (5.37). The recurrence relations for the coefficients

$$
\begin{gather*}
{\left[i(2 k+1)-\left(\Delta+\delta^{\prime} n\right)\right] c_{n, k}=\sqrt{2 \delta^{\prime}} c_{0, k}} \\
2 c_{0, k}=\sqrt{2 \delta^{\prime}} \sum_{n=1}^{\infty} c_{n, k-1} \tag{5.172}
\end{gather*}
$$

yield

$$
\begin{align*}
c_{0, k} & =c_{0, k-1} \sum_{n=1}^{\infty} \frac{\delta^{\prime}}{\left[i(2 k-1)-\left(\Delta+\delta^{\prime} n\right)\right]}  \tag{5.173}\\
& =c_{0, k-1} \int_{0}^{X \rightarrow \infty} \frac{d x}{[i(2 k-1)-\Delta-x]}=c_{0, k-1} \ln \frac{i(2 k-1)-\Delta}{-X}
\end{align*}
$$

and result in

$$
\begin{equation*}
c_{0, k}=\prod_{p=1}^{k} \ln \frac{2 p-1+i \Delta}{i X} \tag{5.174}
\end{equation*}
$$

such that

$$
\begin{equation*}
\psi_{0}(\theta)=\sum_{k=0}^{\infty} \frac{\left(-i \theta^{2}\right)^{k}}{k!} \prod_{p=1}^{k} \ln \frac{2 p-1+i \Delta}{i X} \tag{5.175}
\end{equation*}
$$

In the calculation of the sum (5.175), we make use of the fact that the bandwidth $X$ is large. We employ the representation

$$
\begin{equation*}
c_{0, k}=\exp \left\{\sum_{p=1}^{k} \ln \left[\ln \frac{2 p-1+i \Delta}{i X}\right]\right\} \tag{5.176}
\end{equation*}
$$

and develop the double logarithm in the exponent in the Taylor series over a small parameter $1 / \ln X$ allowing for the fact that the leading term $\ln [-\ln X]$ of this series is non analytic at $1 / \ln X \rightarrow 0$. Although in the first-order approximation the coefficients

$$
\begin{equation*}
c_{0, k}=k \exp \ln [-\ln X]-\frac{1}{\ln X} \sum_{p=1}^{k} \ln \left(\frac{2 p-1}{i}+\Delta\right) \tag{5.177}
\end{equation*}
$$



Fig. 5.14. Integration contour $C$ circumventing the negative integer points in the complex plane can be deformed as $C \rightarrow C_{p}+C^{\prime}$ such that it passes the saddle point $S$. The loop $C^{\prime}$ which goes around the branching $(B)$ point $y=(1+i \Delta) / 2$ allows for the dying oscillations of the populations, whereas the principal part $C_{p}$ of the contour gives the main contribution.
can be found explicitly

$$
\begin{equation*}
\sum_{p=1}^{k} \ln \left(\frac{2 p-1}{i}+\Delta\right)=\ln \prod_{p=1}^{k}\left(\frac{2 p-1}{i}+\Delta\right)=\ln \frac{2^{k}}{i^{k}} \frac{\Gamma\left(k+\frac{1+i \Delta}{2}\right)}{\Gamma\left(\frac{1+i \Delta}{2}\right)} \tag{5.178}
\end{equation*}
$$

in order to be consistent within the approximation one has to take the asymptotic $k, \Delta \gg 1$ for the $\Gamma$-functions and substitute $\ln \Gamma(z)=(z-1 / 2) \ln z-$ $z+\ln \sqrt{2 \pi}$. Equation (5.175) then yields

$$
\begin{equation*}
\psi_{0}(\theta)=\sum_{k=0}^{\infty} \frac{\left(i \theta^{2} \ln X\right)^{k}}{k!} \frac{\left(\frac{\Delta-i-2 i k}{2}\right)^{-(2 k+i \Delta) / 2 \ln X}}{\left(\frac{\Delta-i}{2}\right)^{-i \Delta / 2 \ln X}} \tag{5.179}
\end{equation*}
$$

For small $\theta$, equations (5.175) and (5.179) are both equally good for the numerical evaluation of the probability amplitude of the initial state. However for large $\theta$, the series starts to explode and become inconvenient for numerical calculations. There is an alternative integral representation

$$
\begin{equation*}
\psi_{0}(\theta)=\int_{C} \frac{\left(-i \theta^{2} \ln X\right)^{-y}}{2 \pi i} \Gamma(y) \frac{\left(\frac{\Delta-i}{2}+i y\right)^{(2 y-i \Delta) / 2 \ln X}}{\left(\frac{\Delta-i}{2}\right)^{-i \Delta / 2 \ln X}} d y \tag{5.180}
\end{equation*}
$$



Fig. 5.15. A level near the continuum edge with an exponentially rising coupling $V=V_{0} e^{\alpha t}$. Population of the level (a) as a function of the integrated interaction $\theta^{2}=\int d t V_{0}^{2} e^{2 \alpha t} / \beta=W e^{2 \alpha t} / 2 \alpha$, and the detuning $\Delta_{0}$ of the level from the lower edge of the band. (b) Probability of the transition as a function of the growthrate $\alpha$ for the case of phase decorrelated double interaction: the exponential increase followed by an exponential decrease of the same rate. Small detuning $\Delta_{0} \sim W \sim$ 0.001 (solid line) and large detuning $\Delta_{0} \sim 20 W$ (dashed line). The maximum of the solid curve corresponds to the velocity $\alpha \sim \pi W \ln (X / W) / 2$, where the band width $X$ is taken to be unity.
which relies on the fact that the $\Gamma$-function has simple poles with the residuals $(-1)^{k} / k$ ! all at negative integer points $-k$, that the integration contour $C$ shown in Fig. 5.14 has to circumvent.

For large $\theta$ one can make use of the asymptotic expression $\Gamma(y)=$ $y^{y} e^{-y} / \sqrt{2 \pi y}$ for the $\Gamma$-function

$$
\begin{equation*}
\psi_{0}(\theta)=\int_{C} \frac{\left(\frac{\Delta-i}{2}\right)^{i \Delta / 2 \ln X}}{\sqrt{2 \pi y}\left(\frac{\Delta-i}{2}+i y\right)^{(i \Delta-2 y) / 2 \ln X}}\left(\frac{i y}{e \theta^{2} \ln X}\right)^{y} d y \tag{5.181}
\end{equation*}
$$

and move the integration contour at the saddle point located near $y=$ $-i \theta^{2} \ln X$. For moderate $\theta^{2}$ the integral can be evaluated numerically with the help of the saddle point method, which in this case requires numerical solution of the equation determining the position of the saddle point, and in the extreme of large $\theta^{2}(5.181)$ yields for the population

$$
\begin{equation*}
\left|\psi_{0}(\theta)\right|^{2} \simeq \exp \left[-\frac{4 \theta^{2}+\Delta / \ln X}{2 \theta^{2} \ln X+\Delta}+\frac{\Delta}{\ln X} \arctan \frac{1}{\Delta}\right] \tag{5.182}
\end{equation*}
$$

In Fig. 5.15(a) we show the population $\rho_{t r}=1-\left|\psi_{0}(\theta)\right|^{2}$ transferred from the discrete level as a function of the integrated interaction $\theta^{2}=\int^{t} d t^{\prime} V^{2}\left(t^{\prime}\right) / \beta=$
$\int^{t} d t^{\prime} W\left(t^{\prime}\right)$ and the detuning $\Delta_{0}=\Delta 2 \alpha$. At detunings smaller than the growthrate $\alpha$ of the interaction, one observes dying oscillations, typical of the coherent damping discussed earlier in Chap. 4. It is worth mentioning that the oscillation frequency amounts to $\ln X$, which indicates the important role of the cooperative behavior which was previously discussed on p. 76. For detunings much larger than the growthrate, the transferred population remains zero as long as the interaction remains smaller than the detuning. Note an interesting phenomenon: even for small detunings the transition probability remains small, being of the order of $2 / \ln X$. This somewhat strange effect becomes clear when we recall that the repulsion of the level from the band resulting by analogy to (5.159) in the adiabatic energy

$$
\begin{equation*}
E_{a d}(t)-\Delta_{0}=W \ln \left(\frac{E_{a d}(t)}{-X}\right) \tag{5.183}
\end{equation*}
$$

differs from the probability $W$ of the resonant decay approximately by the factor $\ln X$, which means that even for small detunings, the adiabatic state remains separated from the bottom of the continuum by a gap broader than the width of the Lorentzian profile (3.36) of the resonant transition.

## 6 Time-Dependent Complex Systems

We now turn to time-dependent multilevel systems, which are similar to the two-band models discussed in Chap. 4. The main difference between them and the time-dependent Level-band problem discussed in Chap. 5 is in the topology of the transitions; for two-band models the transition topology is much more involved, and a transition from one level to another does not necessarily imply the immediate inverse transition. On the contrary, the next transition rather occurs to a different energy eigenstate of the unperturbed Hamiltonian. The algebraic structure of these problems is such that they cannot be reduced to simpler systems, and therefore no explicit solution is possible without taking advantage of the statistical models. Therefore the main technique employed for the description is based on the ensemble average over the matrix elements of the coupling, and on the topological selection rules for leading terms in each order of the perturbation series, followed by the exact analytic summation of these series, in complete analogy with the approach developed in Chap. 4.

We start with the problem of a degenerate level crossing a band, resembling the Demkov-Osherov case. However, the presence of the many quantum states in the degenerate level drastically changes the structure of the solution, such that it can no longer be exactly solved by the Laplace contour integral method. The moving level becomes split to individual components and does not lose its population exponentially, but retrieves it when the new component passes the same resonance with a band state that has been crossed by the former component.

We continue by considering perturbations of multilevel systems by a time-dependent random matrix. A generic time-dependent perturbation contains different harmonics, which provoke transitions among the eigenstates of the unperturbed Hamiltonian separated by energies equal to the energies of quanta corresponding to the harmonics. Therefore, under the action of a time-dependent perturbation, the energy distribution can differ dramatically from the Lorentzian distribution typical of the time-independent perturbations switched on at a given moment of time. The ensemble averaged population distribution can be found exactly for a quntum system perturbed by a random matrix with an amplitude arbitrarily changing in time. We consider this problem prior to calculation of the responses of the systems to
the random perturbations. Note that a non-trivial coherent response exists only in the case of the perturbations that belong to the Gaussian orthogonal ensemble.

We also consider the special cases of harmonic time-dependent perturbations for the case of an abrupt and adiabatic switching. These results can be generalized to the bi-harmonic time dependence and to the case of a two band system. It turns out that the energy distribution of the populations over the bands is sensitive to the photon statistics in the case when the timedependent perturbation is considered as a quantized electromagnetic field.

We conclude the chapter by investigating the possibility to control the dynamics of complex quantum systems. The commutation relations among the unperturbed Hamiltonian, the perturbation, and their commutators that make the dynamics complex, also remove all of the holonomic constraints protecting the system from complete control. For the example of the Bloch vectors for the simplest two-level system, we illustrate the geometrical meaning of the algorithms that allow one to determine the time dependence of the perturbation, which yields a desired time evolution of the system. We also show the geometrical meaning of the constraints that exist when the commutation relations correspond to a simple, small-dimensional algebra. We illustrate how control of the complex system can protect a part of it from the loss of coherence.

### 6.1 Degenerate Level Crosses an Infinite Band

Our initial system consists of a level with linearly increasing energy, which crosses an infinite band. In many respects this problem is similar to the Demkov-Osherov problem considered in Sect.5.4.1. The two main differences are in the fact that the level is degenerate, as was the case in Sect. 4.4 for a similar time independent system, and in that the band is infinite on both sides. We also assume that initially, at $t=0$, only one sublevel of the degenerate level is populated, which from the physical point of view implies that the population has been injected into this sublevel while the degenerate level was already in resonance with the band.

The corresponding Schrödinger equation then reads

$$
\begin{align*}
& i \dot{\psi}_{n}=n \delta \psi_{n}+\sum_{m=1}^{N} V_{n m} \psi_{m} \\
& i \dot{\psi}_{m}=\alpha t \psi_{m}+\sum_{n=-\infty}^{\infty} V_{m n} \psi_{n} \tag{6.1}
\end{align*}
$$

and implies that an $N$-fold degenerate level of energy $\alpha t$ moving linearly in time, enters across an infinite band of levels spaced by small energy intervals
$\delta$, which we consider vanishing $\delta \rightarrow 0$. The index $n$ enumerates the levels of the band, and $m$ denotes the sublevels of the degenerate level.

We take the position of the degenerate level at $t=0$ as the energy reference point and attribute the detuning $\alpha t$, linearly depending on time, to the levels $n$ of the band. After such a replacement, Fourier-Laplace transformation of (6.1) yields

$$
\begin{align*}
& \varepsilon \psi_{n}(\varepsilon)=n \delta \psi_{n}(\varepsilon)-i \alpha \frac{\partial \psi_{n}(\varepsilon)}{\partial \varepsilon}+\sum_{m=1}^{N} V_{n m} \psi_{m}(\varepsilon) \\
& \varepsilon \psi_{m}(\varepsilon)=i \delta_{m, 1}+\sum_{n=1}^{\infty} V_{m n} \psi_{n} \tag{6.2}
\end{align*}
$$

where $\delta_{m, m^{\prime}}$ is the Kronecker $\delta$. Equation (6.2) differs from (5.138) not only in the fact that we allow here for the sublevels $m$ of the degenerate level and move the time dependence of the degenerate level detuning to the states of the band, but also in the presence of the initial condition $i \delta_{m, 1}$, which means that initially only the first sublevel of the degenerate level is populated.

Let us now compare (6.2) with (4.5) for the two-band system of Sect.4.2.1. at $\Delta_{m}=0$. The main discrepancy is in the fact that the detunings of the band states $\Delta_{n}=\delta n-i \alpha \frac{\partial}{\partial \varepsilon}$ include the derivative over the spectral variable $\varepsilon$, and therefore after the ensemble average, the renormalized Greens functions $X_{n}(\varepsilon), X_{m}(\varepsilon)$ given by (4.20), (4.21) and corresponding to the static degenerate-level-band problem of (4.5) became, for the time dependent system of (6.2) considered, integral operators $\widehat{X}_{n}, \widehat{X}_{m}$ over $\varepsilon$, that read

$$
\begin{align*}
& \widehat{X}_{n}=\frac{1}{\varepsilon-\left(n \delta-i \alpha \frac{\partial}{\partial \varepsilon}\right)-\left\langle V^{2}\right\rangle N \hat{X}_{m}} \\
& \widehat{X}_{m}=\frac{1}{\varepsilon-\left\langle V^{2}\right\rangle \sum_{n} \widehat{X}_{n}} \tag{6.3}
\end{align*}
$$

In complete analogy with Sect.4.4.1 we replace the sum by the integral and obtain again (4.71)-(4.73), that is

$$
\begin{align*}
\sum_{n} \widehat{X}_{n} & =i \pi g \\
\widehat{X}_{m} & =X(\varepsilon)=\frac{1}{\varepsilon+i \pi g\left\langle V^{2}\right\rangle} \tag{6.4}
\end{align*}
$$

regardless of the presence of the derivative in the denominator of the first of (6.3). Here $g$, as earlier, is the density of the band states. This yields

$$
\begin{equation*}
\widehat{X}_{n}(\varepsilon)=\frac{1}{\varepsilon-\left(n \delta-i \alpha \frac{\partial}{\partial \varepsilon}\right)-\left\langle V^{2}\right\rangle N\left(\varepsilon+i \pi g\left\langle V^{2}\right\rangle\right)^{-1}} \tag{6.5}
\end{equation*}
$$

and by analogy for the complex conjugate Green's function

$$
\begin{align*}
\widehat{X}_{n}(\xi) & =\frac{1}{\xi-\left(n \delta-i \alpha \frac{\partial}{\partial \xi}\right)-\left\langle V^{2}\right\rangle N\left(\xi-i \pi g\left\langle V^{2}\right\rangle\right)^{-1}} \\
X(\xi) & =\frac{1}{\xi-i \pi g\left\langle V^{2}\right\rangle} . \tag{6.6}
\end{align*}
$$

We also note that apart from the replacements $n \leftrightarrow m$ and $2 \leftrightarrow 0$ expressions (4.27) remain valid for the total population of the degenerate level $\rho_{0}(\varepsilon, \xi)$ and of the band $\rho_{1}(\varepsilon, \xi)$, and therefore after having evaluated the sum

$$
\begin{equation*}
\sum_{n} \widehat{X}_{n}(\varepsilon) \widehat{X}_{n}(\xi)=\frac{-2 \pi i g}{\varepsilon-\xi+\left\langle V^{2}\right\rangle N(X(\xi)-X(\varepsilon))+i \alpha\left(\frac{\partial}{\partial \varepsilon}+\frac{\partial}{\partial \xi}\right)} \tag{6.7}
\end{equation*}
$$

by calculating the correspondent integral, one finds for the total population

$$
\begin{align*}
\rho_{1} & =\frac{-2 \pi i V^{2} g}{\varepsilon-\xi+\left\langle V^{2}\right\rangle N(X(\xi)-X(\varepsilon))-i \alpha\left(\frac{\partial}{\partial \varepsilon}+\frac{\partial}{\partial \xi}\right)} \rho_{0} \\
\rho_{0} & =V^{2} N X(\varepsilon) X(\xi) \rho_{1}+X(\varepsilon) X(\xi) . \tag{6.8}
\end{align*}
$$

After substitution of $X(\xi)$ and $X(\varepsilon)$ from (6.4) and (6.6) this immediately results in

$$
\begin{equation*}
\left(\varepsilon-\xi+\frac{\mathcal{V}^{2}(\varepsilon-\xi)}{(\varepsilon+i W)(\xi-i W)}-i \alpha\left(\frac{\partial}{\partial \varepsilon}+\frac{\partial}{\partial \xi}\right)\right) \rho_{1}=\frac{-1}{\varepsilon+i W} \frac{2 i W}{\xi-i W} \tag{6.9}
\end{equation*}
$$

where $\left\langle V^{2}\right\rangle N=\mathcal{V}^{2}$ is the squared cooperative coupling (see p. 76) and $\pi g\left\langle V^{2}\right\rangle=W$ is the probability amplitude decay rate (3.34). After the replacement $\varepsilon-\xi \rightarrow 2 v W, \varepsilon+\xi \rightarrow 2 u W$ this yields the equation

$$
\begin{equation*}
\left(2 v+\frac{\mathcal{V}^{2} 2 v / W^{2}}{u^{2}-(v+i)^{2}}+\frac{i \alpha}{W^{2}} \frac{\partial}{\partial u}\right) \rho_{1}=\frac{-2 i / W^{2}}{u^{2}-(v+i)^{2}} \tag{6.10}
\end{equation*}
$$

By standard methods, one finds the solution of the first-order differential equation (6.10) in the form

$$
\begin{equation*}
\rho_{1}=\frac{-2 i}{\alpha}\left(\frac{u-v-i}{u+v+i}\right)^{-b} \int_{-\infty}^{0} \exp \left(\frac{-2 i v \lambda}{\alpha / W^{2}}\right) \frac{(u+\lambda-v-i)^{b-1}}{(u+\lambda+v+i)^{b+1}} d \lambda \tag{6.11}
\end{equation*}
$$

where $b=-\mathcal{V}^{2} v / \alpha(1-i v)$, or in the original notation

$$
\begin{equation*}
\rho_{1}=\frac{W}{\alpha}\left(\frac{\varepsilon+i W}{\xi-i W}\right)^{b} \int_{-\infty}^{0} e^{i \lambda(\xi-\varepsilon) / \alpha} \frac{(\xi+\lambda-i W)^{b-1}}{(\varepsilon+\lambda+i W)^{1+b}} d \lambda \tag{6.12}
\end{equation*}
$$

where $b=\frac{-\mathcal{V}^{2}(\xi-\varepsilon)}{\alpha(2 W+i(\xi-\varepsilon))}$. Another representation of the integral

$$
\begin{equation*}
\rho_{1}=\frac{W}{\alpha} \frac{1}{(\varepsilon+i W)(\xi-i W)} \int_{-\infty}^{0} e^{i \lambda(\xi-\varepsilon) / \alpha} \frac{\left(1+\frac{\lambda}{\xi-i W}\right)^{b-1}}{\left(1+\frac{\lambda}{\varepsilon+i W}\right)^{1+b}} d \lambda \tag{6.13}
\end{equation*}
$$

is more convenient for further calculations and the asymptotic analysis.
Let us first consider the limit $\alpha \rightarrow 0$. The main contribution to the integral comes from the domain of small $\lambda$ which allows one to approximate (6.13) and write

$$
\begin{align*}
\rho_{1} & =\frac{W}{\alpha} \frac{1}{(\varepsilon+i W)(\xi-i W)} \int_{-\infty}^{0} e^{i \lambda(\xi-\varepsilon) / \alpha+\frac{\lambda}{\xi+i W}(b-1)-\frac{\lambda}{\varepsilon-i W}(b+1)} d \lambda \\
& =\frac{W / \alpha}{(\varepsilon+i W)(\xi-i W) i(\xi-\varepsilon) / \alpha+(\varepsilon+i W)(b-1)-(\xi-i W)(b+1)} . \tag{6.14}
\end{align*}
$$

After substitution of $b$, replacements $(\varepsilon+\xi) / 2 \rightarrow \eta,(\varepsilon-\xi) \rightarrow \zeta$ and integration over $\eta$ this yields

$$
\begin{equation*}
\rho_{1}=\frac{2 \pi i W}{\sqrt{4 \alpha^{2}+\left(4 V^{2}-(2 i W+\zeta)^{2}\right) \zeta^{2}}} \tag{6.15}
\end{equation*}
$$

which for $\alpha=0$ after inverse Fourier transformation over $\zeta$ results in (4.80) for the time-independent system of a band and a degenerate level. This equation shows that the ratio of band population to the degenerate level population in this extreme is a constant, in accordance with (4.82).

Since the limit just considered implies $\alpha T \rightarrow 0$, that is $\alpha /(\varepsilon-\xi) \rightarrow 0$, it cannot be employed in the physically interesting case of $\alpha \rightarrow 0, \alpha T=$ const corresponding to the adiabatic displacement of the degenerate level by a finite energy interval $\alpha T$, which implies that $(\varepsilon-\xi) \sim \alpha$. In order to describe the dynamics of populations in this limit, one has to make the replacements $(\varepsilon+\xi) / 2 \rightarrow \eta,(\varepsilon-\xi) \rightarrow \zeta$ and $-\mathcal{V}^{2} \zeta / \alpha(i \zeta-2 W) \rightarrow b$ in (6.13) and find

$$
\begin{equation*}
\rho_{1}=\frac{W}{\alpha}\left(\frac{\eta+\zeta / 2+i W}{\eta-\zeta / 2-i W}\right)^{b} \int_{-\infty}^{0} e^{i \lambda \zeta / \alpha} \frac{(\eta-\zeta / 2+\lambda-i W)^{b-1}}{(\eta+\zeta / 2+\lambda+i W)^{1+b}} d \lambda . \tag{6.16}
\end{equation*}
$$

Now we note that the integration to be performed over $d \eta$ from $-\infty$ to $\infty$ along the real axis is equivalent to the integration along a contour circumventing two singularities in the upper part of the complex plane at the points $\eta=\zeta / 2-\lambda+i W$ and $\eta=\zeta / 2+i W$. This can be done explicitly after the replacement

$$
\begin{equation*}
\eta \rightarrow-\frac{(2 i W+\zeta)(2 i W+\zeta+\lambda+\lambda x)}{2(2 i W+\zeta+\lambda-\lambda x)} \tag{6.17}
\end{equation*}
$$

which reduces to an integral over $d x$ coinciding with the integral representation of the hypergeometric function and yields

$$
\begin{equation*}
\rho_{1}=\int_{-\infty}^{0} \frac{-2 i W}{\alpha(\zeta+2 i W)} e^{i \lambda \zeta / 2 \alpha}\left(\frac{\zeta-\lambda+2 i W}{\zeta+\lambda+2 i W}\right)^{b}{ }_{2} F_{1}\left(b,-b, 1, \frac{\lambda^{2}}{\lambda^{2}-(\zeta+2 i W)^{2}}\right) d \lambda . \tag{6.18}
\end{equation*}
$$

In this exact expression we can neglect $\zeta$ compared to $W$ in the long-time limit, since the crucial dependence on $\zeta$ already enters the indices $\pm b$ of the hypergeometric function, and after introducing the variables $w=W / \mathcal{V}$, $y=\lambda w, x=\zeta / 2 \alpha w, \Theta=\alpha T / \mathcal{V}$ we obtain the inverse Fourier transform

$$
\begin{align*}
\rho_{1}(\Theta)= & \frac{-1}{\pi} \int_{-\infty}^{0} \exp \left(i x y-2 i \Theta x+\frac{2 i x}{w} \arctan \left(\frac{y}{w}\right)\right) \\
& { }_{2} F_{1}\left(\frac{x}{w},-\frac{x}{w}, 1, \frac{y^{2}}{w^{2}+y^{2}}\right) d x d y . \tag{6.19}
\end{align*}
$$

Now one has to take into account that the Fourier transformation of this hypergeometric function over the indices reads

$$
\begin{equation*}
\int_{-\infty}^{\infty}{ }_{2} F_{1}(u,-u, 1, z) \cos (u v) d u=\frac{\cos \left(\frac{v}{2}\right)}{\sqrt{z-\sin ^{2}\left(\frac{v}{2}\right)}} \tag{6.20}
\end{equation*}
$$

which can be found by a certain exponential replacement in the corresponding integral representation of ${ }_{2} F_{1}$. The spectrum of the hypergeometric function ranges from $v=-2 \arcsin \sqrt{z}$ to $v=2 \arcsin \sqrt{z}$, and therefore after the replacement $y \rightarrow y / w$ one arrives at

$$
\begin{align*}
\rho_{1}(\Theta) & =\int_{-\infty}^{0} \frac{\cos \left(\frac{y}{2}-\Theta w+\arctan \frac{y}{w^{2}}\right)}{\sqrt{\sin ^{2}\left(\arctan \frac{y}{w^{2}}\right)-\sin ^{2}\left(\frac{y}{2}-\Theta w \arctan \frac{y}{w^{2}}\right)}} \frac{d y}{\pi} \\
& =\int_{-\infty}^{0} \frac{\cos \left(\frac{y}{2}-\Theta w+\arctan \frac{y}{w^{2}}\right)}{\sqrt{\sin \left(\Theta w-\frac{y}{2}+2 \arctan \frac{y}{w^{2}}\right) \sin \left(\frac{y}{2}-\Theta w\right)}} \frac{d y}{\pi}, \tag{6.21}
\end{align*}
$$

where the arguments of sine and cosine must remain within the interval $(-\pi, \pi)$. Here we have also taken into account, that $\sin ^{2}(\arctan z)=$ $z^{2} /\left(1+z^{2}\right)$.

In Fig. 6.1 we show the population $\rho_{0}(\Theta)=1-\rho_{1}(\Theta)$ of the degenerate level as a function of the displacement $\Delta=\alpha T$ and the coupling parameter $\mathcal{V} / W=1 / w$, calculated with the help of (6.21). Note that $1 / w^{2}$ equates to the ratio of the number of states $N$ comprising the degenerate level and the number $W g$ of the band levels in resonance, discussed on p. 76. One sees the main trend of the population behavior, namely an almost linear decrease of $\rho_{0}(T)$ toward zero with a slope proportional to the parameter $w$. This is a natural consequence of the fact that in the course of the motion of the degenerate level across the band, the degeneracy is raised by the interaction. Therefore the populations of the resulting components make transitions to the states of the band via the mechanism of adiabatic population transport. Each transition can occur in both directions: from a state of the level to a state of the band and vice versa, depending on which one of these two states


Fig. 6.1. Population (a) of the degenerate level as a function of the dimensionless energy shift $\Theta=\Delta W / V^{2} N$ and the scaled coupling $1 / w=\sqrt{V^{2} N / W^{2}}$ in the limit $\alpha \rightarrow 0$, and $T \alpha \rightarrow \Delta$, for the system of an $N$-fold degenerate level interacting with a band. The eigenvalues of the system (b) are shown as a function of time. In order to distinguish among the eigenstates originating from the degenerate level and those originating from the states of the band we took a finite time for the interaction switch. The density of states that belong to the degenerate level are shown by dashed lines for two different times. The population of the adiabatic levels (dashdot line) given by the initial distribution of the degenerate level's density of states propagated along the adiabatic curves after convolution with the final distribution of these states gives the population (a).
corresponds to the populated adiabatic level before the crossing. The number of such transitions is proportional to the product of the band state density and the displacement, whereas the population of each component is inversely proportional to the degeneracy $N$ of the level. One can see this phenomenon by tracing the time evolution of the eigenstates shown in Fig. 6.1b). The population of the adiabatic states does not change, whereas the population of the degenerate level at a given moment in time depends on the number of the initially populated adiabatic states that still belong to this manifold at that moment.

On a long-time scale, the band population probability approaches unity. Deviation from unity can be found from (6.21) when we denote $\frac{y}{2}-\Theta w+$ $\arctan \frac{y}{w^{2}}=u$, take this quantity as the integration variable, and find approximately the derivative $\partial u / \partial y \simeq 1 / 2+1 /\left(w^{2}+4 \Theta^{2}\right)$. This yields

$$
\begin{align*}
\rho_{1}(\Theta) & =\int \frac{d u}{\sqrt{\sin ^{2}\left(\arctan \frac{2 \Theta w}{w^{2}}\right)-u^{2}}} \frac{1}{\pi \frac{\partial u}{\partial y}} \\
& \simeq \frac{1}{1+2 /\left(w^{2}+4 \Theta^{2}\right)} \simeq 1-\frac{\mathcal{V}^{2}}{2(\alpha T)^{2}} \tag{6.22}
\end{align*}
$$

that is

$$
\begin{equation*}
\rho_{0}(T) \simeq \frac{\mathcal{V}^{2}}{2(\alpha T)^{2}} \tag{6.23}
\end{equation*}
$$

In this expression one recognizes the population $\rho_{0}$ of (3.19) for a single level separated from the resonance with another populated single level by a large detuning $\alpha T$ and coupled to this level by an interaction $\mathcal{V} / 2$, which can be attributed to the joint effect of the wings of the initial population distribution and the influence of the second-order perturbation of the level's states by the populated states of the band.

An explicit analytical solution for another important limit of high velocities, that is $\alpha \rightarrow \infty$, can be easily found after the replacement $\lambda \rightarrow \alpha \lambda$ in (6.18) which yields, in the limit under consideration

$$
\begin{align*}
\rho_{1} & =\int_{-\infty}^{0} \frac{-2 i W e^{i \lambda \zeta / 2}}{(\zeta+2 i W)}(-1)^{b}{ }_{2} F_{1}(b,-b, 1,1) d \lambda \\
& =\frac{-i 4 W}{2 \pi b(\zeta+2 i W) \zeta}\left(1-e^{2 i \pi b}\right) \tag{6.24}
\end{align*}
$$

where the relation

$$
\begin{equation*}
(-1)^{b}{ }_{2} F_{1}(b,-b, 1,1)=\frac{i}{2 \pi b}\left(1-e^{2 i \pi b}\right) \tag{6.25}
\end{equation*}
$$

has been employed. We have also to set $b \rightarrow 0$, and hence

$$
\begin{equation*}
\rho_{1}=\frac{-4 W}{(\zeta+2 i W) \zeta} \tag{6.26}
\end{equation*}
$$

which after inverse Fourier transformation gives a transparent expression

$$
\begin{equation*}
\rho_{1}(t)=1-e^{-2 \pi g\left\langle V^{2}\right\rangle t} \tag{6.27}
\end{equation*}
$$

in complete analogy with (3.34). This means that, at such a high velocity of displacement of the degenerate level across the band, each component of the level encounters such a large number of band states that it decays to these states individually, as if they were a real continuum, while all of the interference phenomena such as recurrences and revivals are completely suppressed.

### 6.2 Perturbation Proportional to a Random Matrix

Thus far, we have considered particular forms of the time dependencies of the interaction or the level positions, that allow one to find the result analytically. Looking at this class of problems from the viewpoint of the algebraic properties of quantum systems discussed in Sect. 5.1, one can say that their algebraic structure is such that the representation

$$
\begin{equation*}
\widehat{U}(t)=e^{\sum_{n=0}^{N_{O}} f_{n}(t) \widehat{\mathcal{G}}_{n}} \tag{6.28}
\end{equation*}
$$

of the evolution operator $\widehat{U}(t)$ given in terms of the generators $\widehat{\mathcal{G}}_{n}$ of the unitary group, by analogy to (5.5), has a very special relationship among the functions $f_{n}(t)$, that allow one to express the solution in terms of integrals or special functions, something that is completely impossible for a generic problem. In this section we turn to the problem of a perturbation with an arbitrary time dependence, which does not have such a hidden symmetry. In order to achieve an analytical result for this case, we have to make some simplifying assumptions compensating for the complex character of the exact dynamics. We therefore assume that the Hamiltonian has the form

$$
\begin{equation*}
\widehat{H}(t)=\widehat{H}_{0}+q(t) \widehat{V} \tag{6.29}
\end{equation*}
$$

of an unperturbed Hamiltonian $\widehat{H}_{0}$ with a dense and uniform spectrum, plus a perturbation in the form of a time-independent random matrix $\widehat{V}$ multiplied by an arbitrary time-dependent coefficient $q(t)$. In the representation of the eigenfunctions of $\widehat{H}_{0}$ shown in Fig. 6.2(a) the corresponding Schrödinger equation reads

$$
\begin{equation*}
i \dot{\psi}_{n}=\Delta_{n} \psi_{n}+q(t) \sum_{m} V_{n m} \psi_{m} \tag{6.30}
\end{equation*}
$$

where $\Delta_{n}$ stand for the eigenstate energies.


Fig. 6.2. Model of a collision. Dense spectrum (a) subjected to the action of a perturbation which has the structure $q(t) \widehat{V}$ of a random matrix $\widehat{V}$ multiplied by a time-dependent function $q(t)$. Typical population distribution after the collision (b). Population distribution (c) as a function of the interaction strength $W=\pi g V^{2}$ and the detuning $\Delta$ for a Gaussian $q(t)=e^{-\alpha^{2} t^{2}}$ time dependence.

We now make use of the idea already employed in Sect. 5.1, and replace the Schrödinger equation (6.30) with time dependent coefficients by an equation of a higher dimensionality with time independent coefficients by analogy to (5.8). For the evolution operator of this system in a Fourier representation one obtains

$$
\begin{equation*}
\varepsilon \widehat{U}(\varepsilon, \tau)-\left[\widehat{H}_{0}+q(\tau) \widehat{V}-i \frac{\partial}{\partial \tau}\right] \widehat{U}(\varepsilon, \tau)=\widehat{I} \tag{6.31}
\end{equation*}
$$

that is

$$
\begin{equation*}
\widehat{U}(\varepsilon, \tau)=\left(\varepsilon-\widehat{H}_{0}+i \frac{\partial}{\partial \tau}-q(\tau) \widehat{V}\right)^{-1} . \tag{6.32}
\end{equation*}
$$

In other words

$$
\begin{equation*}
\psi_{n}(t)=\left.\int \frac{e^{-i \varepsilon t}}{2 \pi} \sum_{m}\left(\frac{1}{\varepsilon-\widehat{H}_{0}+i \frac{\partial}{\partial \tau}-q(\tau) \widehat{V}}\right)_{n, m}\right|_{\tau=t} \psi_{m}(t=0) \tag{6.33}
\end{equation*}
$$

This equation takes exactly the same form as (4.10) when we attribute the derivative $-i \frac{\partial}{\partial \tau}$ to the unperturbed Hamiltonian $\widehat{H}_{0}$. Therefore all of the calculations performed in Sect. 4.1 such as the expansion of the evolution operator in a Taylor series over $\widehat{V}$ followed by ensemble averaging can also be repeated for the problem under consideration. The only new requirement is to take into account the order of operators, since $-i \frac{\partial}{\partial \tau}$ and $q(\tau)$ do not commute.

By complete analogy to (4.21), the calculations yield the equation

$$
\begin{equation*}
\widehat{X}_{n}(\varepsilon, \tau)=\frac{1}{\varepsilon-\Delta_{n}+i \frac{\partial}{\partial \tau}-q(\tau)\left\langle V^{2}\right\rangle \sum_{m} \widehat{X}_{m}(\varepsilon) q(\tau)}, \tag{6.34}
\end{equation*}
$$

for the renormalized operators $\widehat{X}_{n}(\varepsilon)$ although now they have to be considered also as operators in $\tau$-space. However for the sum $\widehat{Q}(\varepsilon)=\sum_{m} \widehat{X}_{m}(\varepsilon)$, under the condition that the spectrum of $\widehat{H}_{0}$ is broad enough, we still obtain the solutions (4.32). This yields

$$
\begin{equation*}
\widehat{X}_{n}(\varepsilon, \tau)=\frac{1}{\varepsilon-\Delta_{n}+i \frac{\partial}{\partial \tau}+i \pi g\left\langle V^{2}\right\rangle q^{2}(\tau)} \tag{6.35}
\end{equation*}
$$

### 6.2.1 Population Distribution

We now consider the population $\rho_{n}$ of an arbitrary state $|n\rangle$. Since we have just one band, the perturbation series for the population is a sum of all the diagrams shown in Fig. 6.3 similar to Fig. 4.12, and $\rho_{n}$ is given by (4.25) with the only difference that instead of (4.24) the factor $\Xi$ reads

$$
\begin{align*}
\widehat{\Xi}(\varepsilon, \xi, \tau, \theta) & =\sum_{n} \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) V^{2} q(\tau) q(\theta) \\
& =\int_{-\infty}^{\infty} \frac{g\left\langle V^{2}\right\rangle}{\varepsilon-\Delta+i \frac{\partial}{\partial \tau}+i \pi g V^{2} q^{2}(\tau)} \frac{d \Delta}{\xi-\Delta-i \frac{\partial}{\partial \theta}-i \pi g V^{2} q^{2}(\theta)} q(\tau) q(\theta) \\
& =\frac{2 i \pi g V^{2}}{\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left(q^{2}(\tau)+q^{2}(\theta)\right)} q(\tau) q(\theta) . \tag{6.36}
\end{align*}
$$



Fig. 6.3.

Here the variable $\theta$ plays for $X_{n^{\prime}}(\xi, \theta)$ the same role as $\tau$ does for $X_{n^{\prime}}(\varepsilon, \tau)$, and the integral replacing the sum has been calculated by taking the residual in the pole of the upper part of the complex plane. The sum of the series (4.25) for the initial distribution $\rho_{m}(t=0)$ yields for the resulting population

$$
\begin{gather*}
\rho_{n}(\varepsilon, \xi, \tau, \theta)=\widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \rho_{n}(t=0) \\
+\widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) V^{2} q(\tau) q(\theta) \frac{1}{1-\widehat{\Xi}(\varepsilon, \xi)} \sum_{m} \widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}(\xi, \theta) \rho_{m}(t=0) \tag{6.37}
\end{gather*}
$$

We note that the relation

$$
\begin{equation*}
\frac{1}{1-\widehat{\Xi}(\varepsilon, \xi)}=\frac{\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left(q^{2}(\tau)+q^{2}(\theta)\right)}{\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}(q(\tau)-q(\theta))^{2}} \tag{6.38}
\end{equation*}
$$

follows from (6.36), and consider the initial condition $\rho_{m}(t=0)=\delta_{m}^{0}$, where $\Delta_{0}=0$. For $n \neq 0$ (6.37) takes the form

$$
\begin{gather*}
\rho_{n}(\varepsilon, \xi, \tau, \theta)=\widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) V^{2} q(\tau) q(\theta) \widehat{Z}(\xi, \theta, \varepsilon, \tau) \\
{\left[\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left(q^{2}(\tau)+q^{2}(\theta)\right)\right] \widehat{X}_{0}(\varepsilon, \tau) \widehat{X}_{0}(\xi, \theta),} \tag{6.39}
\end{gather*}
$$

where

$$
\begin{equation*}
\widehat{Z}(\xi, \theta, \varepsilon, \tau)=\frac{1}{\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}(q(\tau)-q(\theta))^{2}} \tag{6.40}
\end{equation*}
$$

and the term $\widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \delta_{0}^{n}$ is omitted. We now note the consequence of (6.35)

$$
\begin{equation*}
\left[\widehat{X}_{m}^{-1}(\varepsilon, \tau)-\widehat{X}_{m}^{-1}(\xi, \theta)\right]=\left[\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left(q^{2}(\tau)+q^{2}(\theta)\right)\right] \tag{6.41}
\end{equation*}
$$

which is valid for any $m$, including $m=0$. Therefore, the product of the right-hand side of (6.39) with $\widehat{X}_{0}(\varepsilon, \tau) \widehat{X}_{0}(\xi, \theta)$ yields $\left[\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)\right]$, and hence substitution of this product into (6.39) results in

$$
\begin{equation*}
\rho_{n}(\varepsilon, \xi, \tau, \theta)=\widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) V^{2} q(\tau) q(\theta) \widehat{Z}(\xi, \theta, \varepsilon, \tau)\left[\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)\right] \tag{6.42}
\end{equation*}
$$

The relation

$$
\begin{equation*}
\left[\widehat{Z}^{-1}(\xi, \theta, \varepsilon, \tau)+\widehat{X}_{m}^{-1}(\varepsilon, \tau)-\widehat{X}_{m}^{-1}(\xi, \theta)\right]=2 \pi g V^{2} q(\tau) q(\theta) \tag{6.43}
\end{equation*}
$$

allows one to make the last simplification for the general form of operators

$$
\begin{align*}
2 \pi g \rho_{n}(\varepsilon, \xi, \tau, \theta)= & -i \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \\
& {\left[\widehat{Z}^{-1}(\xi, \theta, \varepsilon, \tau)+\widehat{X}_{m}^{-1}(\varepsilon, \tau)-\widehat{X}_{m}^{-1}(\xi, \theta)\right] } \\
& \widehat{Z}(\xi, \theta, \varepsilon, \tau)\left[\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)\right], \tag{6.44}
\end{align*}
$$

which results in

$$
\begin{align*}
& 2 \pi g \rho_{n}(\varepsilon, \xi, \tau, \theta)=-i \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta)\left[\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)\right] \\
& -i\left[\widehat{X}_{n}(\xi, \theta)-\widehat{X}_{n}(\varepsilon, \tau)\right] \widehat{Z}(\xi, \theta, \varepsilon, \tau)\left[\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)\right] . \tag{6.45}
\end{align*}
$$

Further simplifications require the substitution of explicit expressions for $\widehat{X}_{n}(\varepsilon, \tau), \widehat{X}_{n}(\xi, \theta)$, and $\widehat{Z}(\xi, \theta, \varepsilon, \tau)$, that being inverse to the first-order differential operators (6.35), (6.40), are integral operators of the form

$$
\begin{equation*}
\widehat{X}_{n}(\varepsilon, \tau) \psi(\tau)=\int_{-\infty}^{\infty} X_{n}\left(\varepsilon, \tau, \tau^{\prime}\right) \psi\left(\tau^{\prime}\right) d \tau^{\prime} \tag{6.46}
\end{equation*}
$$

For the kernels $Z\left(\xi, \theta, \theta^{\prime}, \varepsilon, \tau, \tau^{\prime}\right), X_{n}\left(\varepsilon, \tau, \tau^{\prime}\right)$, and $X_{n}\left(\xi, \theta, \theta^{\prime}\right)$ of these operators one finds the first-order differential equations

$$
\begin{align*}
{\left[\varepsilon-\xi+i \pi g V^{2}(q(\tau)-q(\theta))^{2}\right] Z\left(\xi, \theta, \theta^{\prime}, \varepsilon, \tau, \tau^{\prime}\right) } & \\
+i\left(\frac{\partial}{\partial \theta}+\frac{\partial}{\partial \tau}\right) Z\left(\xi, \theta, \theta^{\prime} \varepsilon, \tau, \tau^{\prime}\right) & =\delta\left(\tau-\tau^{\prime}\right) \delta\left(\theta-\theta^{\prime}\right), \\
\left(\varepsilon-\Delta_{n}+i \pi g\left\langle V^{2}\right\rangle q^{2}(\tau)\right) X_{n}\left(\varepsilon, \tau, \tau^{\prime}\right)+i \frac{\partial X_{n}\left(\varepsilon, \tau, \tau^{\prime}\right)}{\partial \tau} & =\delta\left(\tau-\tau^{\prime}\right), \\
\left(\xi-\Delta_{n}-i \pi g\left\langle V^{2}\right\rangle q^{2}(\theta)\right) X_{n}\left(\xi, \theta, \theta^{\prime}\right)-i \frac{\partial X_{n}\left(\xi, \theta, \theta^{\prime}\right)}{\partial \theta} & =\delta\left(\theta-\theta^{\prime}\right) . \tag{6.47}
\end{align*}
$$

The solutions

$$
\begin{align*}
X_{n}\left(\varepsilon, \tau, \tau^{\prime}\right) & =-i \Theta\left(\tau-\tau^{\prime}\right) e^{i\left(\varepsilon-\Delta_{n}\right)\left(\tau-\tau^{\prime}\right)-\pi g V^{2}} \int_{\tau^{\prime}}^{\tau} q^{2}(x) d x \\
X_{n}\left(\xi, \theta, \theta^{\prime}\right) & =i \Theta\left(\theta-\theta^{\prime}\right) e^{-i\left(\xi-\Delta_{n}\right)\left(\theta-\theta^{\prime}\right)-\pi g V^{2} \int_{\theta^{\prime}}^{\theta} q^{2}(x) d x} \\
Z\left(\xi, \theta, \theta^{\prime}, \varepsilon, \tau, \tau^{\prime}\right) & =-i \Theta\left(\tau+\theta-\tau^{\prime}-\theta^{\prime}\right) \delta\left(\frac{\tau-\theta-\tau^{\prime}+\theta^{\prime}}{2}\right) \exp [i(\varepsilon-\xi) \\
\times \frac{\tau+\theta-\tau^{\prime}-\theta^{\prime}}{2} & \left.-\pi g V^{2} \int_{\frac{\tau+\theta}{2}}^{\frac{\tau+\theta}{2}}\left(q\left(\frac{\tau-\theta}{2}+x\right)-q\left(x-\frac{\tau-\theta}{2}\right)\right)^{2} d x\right] \tag{6.48}
\end{align*}
$$

substituted into (6.45), integrated over all intermediate indexes, followed by the substitution $\tau \rightarrow t, \theta \rightarrow t$ after the inverse Fourier transformation yield for the spectral population density at $t>0$

$$
\begin{equation*}
\rho_{n}(t)=\operatorname{Re} \int_{0}^{t} \exp \left[i \Delta_{n} v-2 \pi g V^{2}\left(\int_{0}^{t} q^{2}(x) d x-\int_{0}^{t-v} q(x) q(v+x) d x\right)\right] \frac{d v}{\pi g} \tag{6.49}
\end{equation*}
$$

We note that the contribution comes only from the second term of the righthand side of (6.45), since the first term is a $\delta\left(\Delta_{n}\right)$-like one.

We present more details of the calculations. One can obtain the result (6.49) by considering two of six terms of the expanded (6.45), namely $\widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \widehat{X}_{0}(\xi, \theta)$ and $\widehat{X}_{n}(\xi, \theta) \widehat{Z}(\xi, \theta, \varepsilon, \tau) \widehat{X}_{0}(\varepsilon, \tau)$. Two other terms given by the replacement $\varepsilon \rightleftharpoons \xi$ are just the complex conjugate of these terms, and the two remaining terms $\widehat{X}_{n}(\xi, \theta) \widehat{Z}(\xi, \theta, \varepsilon, \tau) \widehat{X}_{0}(\xi, \theta)$ and its complex conjugate vanish, since all of the singularities of this expression are located at the lower (upper) part of the complex plane $\xi$ (upper part of $\varepsilon$ plane for the conjugate). Substitution of (6.48) yields

$$
\begin{align*}
\widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \widehat{X}_{0}(\xi, \theta)= & i \Theta\left(\tau-\tau^{\prime}\right) \int d \theta^{\prime \prime} \Theta\left(\theta-\theta^{\prime \prime}\right) \Theta\left(\theta^{\prime \prime}-\theta^{\prime}\right) \\
& e^{i\left(\varepsilon-\Delta_{n}\right)\left(\tau-\tau^{\prime}\right)-\pi g V^{2}} \int_{\tau^{\prime}}^{\tau} q^{2}(x) d x \\
& e^{-i\left(\xi-\Delta_{n}\right)\left(\theta-\theta^{\prime \prime}\right)-\pi g V^{2}} \int_{\theta^{\prime \prime}}^{\theta} q^{2}(x) d x \\
& e^{-i\left(\xi-\Delta_{n}\right)\left(\theta^{\prime \prime}-\theta^{\prime}\right)-\pi g V^{2} \int_{\theta^{\prime}}^{\theta^{\prime \prime}} q^{2}(x) d x} \tag{6.50}
\end{align*}
$$

which after performing the integration of the inverse Fourier transformation over $\varepsilon$ and $\xi$ results in

$$
\begin{align*}
& \frac{1}{4 \pi^{2}} \int e^{-i(\varepsilon-\xi) t} \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \widehat{X}_{0}(\xi, \theta)= \\
& \Theta(t) \delta\left(\theta^{\prime}\right) \delta\left(\tau^{\prime}\right) e^{-2 \pi g V^{2}} \int_{0}^{t} q^{2}(x) d x  \tag{6.51}\\
& \frac{1-e^{-i \Delta_{n} t}}{i \Delta_{n}}
\end{align*}
$$

The real part of this expression is proportional to $\sin \left(\Delta_{n} t\right) / \Delta_{n}$, which in the long-time asymptotic gives the singular term $\delta\left(\Delta_{n}\right)$, unimportant for nonzero detunings.

The principal term contributing to (6.48) reads

$$
\begin{align*}
& -\widehat{X}_{n}(\xi, \theta) \widehat{Z}(\xi, \theta, \varepsilon, \tau) \widehat{X}_{0}(\varepsilon, \tau)=\Theta\left(\theta-\theta^{\prime \prime}\right) e^{-i\left(\xi-\Delta_{n}\right)\left(\theta-\theta^{\prime \prime}\right)-\pi g V^{2} \int_{\theta^{\prime \prime}}^{\theta} q^{2}(x) d x} \\
& i \Theta\left(\tau+\theta^{\prime \prime}-\tau^{\prime \prime}-\theta^{\prime}\right) \delta\left(\frac{\tau-\theta^{\prime \prime}-\tau^{\prime \prime}+\theta^{\prime}}{2}\right) \Theta\left(\tau^{\prime \prime}-\tau^{\prime}\right) e^{i \varepsilon\left(\tau^{\prime \prime}-\tau^{\prime}\right)-\pi g V^{2}} \int_{\tau^{\prime}}^{\tau^{\prime \prime}} q^{2}(x) d x \\
& e^{i(\varepsilon-\xi) \frac{\tau+\theta^{\prime \prime}-\tau^{\prime \prime}-\theta^{\prime}}{2}-\pi g V^{2} \frac{\int^{\frac{\tau+\theta^{\prime \prime}}{2}}}{\frac{\tau^{\prime \prime}+\theta^{\prime}}{2}}\left(q\left(\frac{\tau+\theta^{\prime \prime}}{2}+x\right)-q\left(\frac{\tau+\theta^{\prime \prime}}{2}-x\right)\right)^{2} d x} \tag{6.52}
\end{align*}
$$

whereas its inverse Fourier transform over $\varepsilon$ and $\xi$ gives the Dirac $\delta$-functions $\delta\left(\frac{\tau-\theta^{\prime \prime}-\tau^{\prime \prime}+\theta^{\prime}}{2}\right)$ and $\delta\left(-t-\tau^{\prime}+\frac{\tau+\theta^{\prime \prime}+\tau^{\prime \prime}-\theta^{\prime}}{2}\right)$ and yields

$$
\begin{align*}
& -\int e^{-i(\varepsilon-\xi) t} \widehat{X}_{n}(\xi, \theta) \widehat{Z}(\xi, \theta, \varepsilon, \tau) \widehat{X}_{0}(\varepsilon, \tau) \frac{d \varepsilon d \xi}{4 \pi^{2}}=i \Theta\left(\theta^{\prime \prime}-\theta^{\prime}\right) \Theta\left(\theta-\theta^{\prime \prime}\right) \\
& \delta\left(t-\theta+\theta^{\prime}\right) e^{i \Delta_{n}\left(\theta-\theta^{\prime \prime}\right)-\pi g V^{2} \int_{\theta^{\prime \prime}}^{\theta} q^{2}(x) d x} \delta\left(\frac{\tau-\theta^{\prime \prime}-\tau^{\prime \prime}+\theta^{\prime}}{2}\right) \delta\left(-t-\tau^{\prime}+\tau\right) \\
& e^{-\pi g V^{2} \frac{\tau+\theta^{\prime \prime}}{2}}\left[q\left(\frac{\tau-\theta^{\prime \prime}}{2}+x\right)-q\left(x-\frac{\tau-\theta^{\prime \prime}}{2}\right)\right]^{2} d x  \tag{6.53}\\
& \frac{\tau+\theta^{\prime}}{2}
\end{align*} \Theta\left(\tau^{\prime \prime}-\tau^{\prime}\right) e^{-\pi g V^{2} \int_{\tau^{\prime}}^{\tau^{\prime \prime}} q^{2}(x) d x},(6.53) .
$$

where we have to set $\theta=\tau=t$. This results in

$$
\begin{align*}
& -\int e^{-i(\varepsilon-\xi) t} \widehat{X}_{n}(\xi, \theta) \widehat{Z}(\xi, \theta, \varepsilon, \tau) \widehat{X}_{0}(\varepsilon, \tau) \frac{d \varepsilon d \xi}{4 \pi^{2}}=-i \Theta\left(t-\tau^{\prime \prime}\right) \Theta\left(\tau^{\prime \prime}\right) \\
& \delta\left(\tau^{\prime}\right) \delta\left(\theta^{\prime}\right) e^{i \Delta_{n} \tau^{\prime \prime}} e^{2 \pi g V^{2}} \int_{0}^{t-\tau^{\prime \prime}} q(x) q\left(x+\tau^{\prime \prime}\right) d x-2 \pi g V^{2} \int_{0}^{t} q^{2}(x) d x \tag{6.54}
\end{align*}
$$

which after the substitution $\tau^{\prime \prime} \rightarrow v$ and after adding the complex conjugate, yields (6.49).

For the interaction vanishing at $|t| \rightarrow \infty$ the integration limits $[0, t]$ have to be extended to $\pm \infty$

$$
\begin{equation*}
\pi g \rho_{n}(t)=\operatorname{Re} \int_{-\infty}^{\infty} \exp \left\{i \Delta_{n} v-2 \pi g V^{2} \int_{-\infty}^{\infty}\left[q^{2}(x)-q(x) q(v+x)\right] d x\right\} d v \tag{6.55}
\end{equation*}
$$

Briefly, the result obtained means that the population distribution over the spectrum of a system with a uniform density of states perturbed by a timedependent random matrix equals the Fourier transform $v \rightarrow \Delta_{n}$ of the exponent of the interaction autocorrelation. Figure 6.2(c) illustrates the situation for a particular example of a gaussian time dependence $q(x)$ resulting in Gaussian autocorrelation.

Note that the result (6.55) can be generalized to the case of timedependent detunings $\Delta_{n}(t)$. To this end the combination $\Delta_{n} v$ has to be replaced by the integral $\int_{-\infty}^{\infty}\left[\Delta_{n}(t+v)-\Delta_{n}(t)\right] d t$, as it follows from the exact solution of the first-order differential equations (6.47). The same generalization can be done for (6.49); in this case the replacing integral reads $\int_{0}^{t} \Delta_{n}(x) d x-\int_{0}^{t-v} \Delta_{n}(x) d x$.

### 6.2.2 Response to Perturbation Proportional to a Random Matrix

Thus far our main concern was the population distribution over states of time-dependent, complex quantum systems. However, it is not the only possible quantity that manifests universal behavior independent of spectrum
details and interaction in the systems. Another example of such a quantity is interaction with external fields. We consider an external field $\mathcal{F}$ where the interaction Hamiltonian $\mathcal{F} \widehat{x}$ is proportional to a generalized coordinate operator $\widehat{x}$. If the time-dependent mean value $x=\langle\Psi(t)| \widehat{x}|\Psi(t)\rangle$ of this operator does not vanish after the ensemble average discussed in Chap. 4, the complex system interacts with such a field, and at the macroscopic level such an interaction is responsible for the polarization of the media corresponding to $\mathcal{F}$. Electric polarization is a typical example, when $\mathcal{F}$ is simply the electric field strength, and $\widehat{x}$ is the dipole moment operator.

## Diagram Series Structure

The terms of perturbation series contributing to this interaction have a similar, although somewhat different structure as compared to the perturbation series for the population shown in Figs. 4.10-4.12. The main difference results from the necessity to find a counterpart for the matrix elements of $\widehat{x}$ among the transition matrix elements $\widehat{V}$ belonging either to the perturbation series for $|\Psi(t)\rangle$ or that for $\langle\Psi(t)|$, such that $\left\langle x_{k m} V_{m k}\right\rangle \neq 0$. Also important is the difference between the statistical ensemble for imaginary and for real matrix elements of the perturbation, concisely written as two possibilities $\left\langle V_{k m} V_{k m}\right\rangle=0$ and $\left\langle V_{k m} V_{k m}\right\rangle=\left\langle V_{k m} V_{m k}\right\rangle=\left\langle V^{2}\right\rangle \neq 0$, respectively. For the Gaussian statistics of the matrix elements these cases correspond respectively to the Gaussian orthogonal ensemble (GOE) and the Gaussian unitary ensemble (GUE). In Fig. 6.4 we show the diagrams that, after the ensemble


Fig. 6.4. Non-vanishing contribution of the perturbation series for the matrix element of a generalized coordinate operator $\hat{x}$. Dotted lines correspond to $|\Psi\rangle$ and dashed lines to $\langle\Psi|$. (a) Case of complex random couplings, typical of the unitary Gaussian ensemble. The main contribution results from the susceptibility $\langle\hat{x} \hat{X}(\varepsilon) \hat{V}\rangle+$ c.c. averaged over the population distribution $\rho$. One recognizes the parallel diagrams contributing to the series for the populations $\rho$, similar to that of Fig. 4.11. (b) Real coupling $\hat{V}$, typical of the orthogonal Gaussian ensemble. Apart from the diagrams for the population one sees part $C$ originating from the averages $\left\langle V_{k m} V_{k m}\right\rangle$. Terms of the same topology are responsible for the Cooper pairs in the theory of superconductivity.
average, give a non-vanishing contribution for the mean matrix element for these two cases.

Thus we have encountered for the first time a quantity for which the results depend on the symmetry with respect to time inversion, a property of random ensembles, typically associated with the presence of a magnetic field. In the absence of magnetic fields, all of the interaction matrix elements are real (or can be selected as real in a certain basis), and therefore $\left\langle V_{k m} V_{k m}\right\rangle=\left\langle V_{k m} V_{m k}\right\rangle$. Therefore, apart from the diagrams where the nonvanishing contribution results either from the transitions coupled to the antiparallel transition in the perturbation series for the same wavefunction or from the transitions coupled with the parallel transition in the perturbation series for the conjugated wavefunctions, we have to take into account the contributions $\widehat{C}$, resulting from the coupling of the antiparallel transitions in series for the different wavefunctions. We note that the diagrams of such a topology are important and well known in the theory of superconductivity, allowing for the formation of Cooper pairs, or "Cooperons". The relative number of these terms is, however, usually smaller by a factor $1 / N$ of the inverse number of levels participating in the process, since topologically an extra self-crossing of the diagram is required for this purpose. That was the reason why these diagrams have not been taken into account for the populations. Nevertheless, for susceptibilities they may give the leading contributions.

Let us consider first the other case, $\left\langle V_{k m} V_{k m}\right\rangle=0$, typical of the systems with no time-reversal symmetry. The principal contribution to the perturbation series for $\langle\Psi(t)| \widehat{x}|\Psi(t)\rangle$ has the structure shown in Fig. 6.4(a) and reads

$$
\begin{align*}
x(t) & =\langle\Psi(t)| \widehat{x}|\Psi(t)\rangle \\
& =\operatorname{Re} \sum_{n, m} \int e^{i(\varepsilon-\xi) t} \widehat{\rho}_{n}(\varepsilon, \xi, \tau, \theta)\left\langle x_{n m} V_{m n}\right\rangle \widehat{X}_{m}(\varepsilon, \tau) \frac{d \varepsilon d \xi}{4 \pi^{2}} \tag{6.56}
\end{align*}
$$

where the operators $\widehat{\rho}_{n}$ and $\widehat{X}_{m}$ are given in (6.35), (6.42) and the effect is proportional to the correlation value $\left\langle x_{n m} V_{m n}\right\rangle$. When this average is independent of the index $m$, the summation over the uniform spectrum yields a purely imaginary value $\sum_{m} \widehat{X}_{m}(\varepsilon, \tau)=\widehat{Q}(\varepsilon)$, such that the result (6.56) vanishes. For the more general case of non-uniform density of states and the correlation $\left\langle x_{n m} V_{m n}\right\rangle$, the combination $\operatorname{Re} \sum_{m}\left\langle x_{n m} V_{m n}\right\rangle \widehat{X}_{m}(\varepsilon, \tau)$ may be interpreted as the linear response of the system in the state $n$, whereas summation over $n$ implies taking the average of the linear response by convolution with the instantaneous population distribution $\widehat{\rho}_{n}(t)=\int e^{i(\varepsilon-\xi) t} \widehat{\rho}_{n}(\varepsilon, \xi, \tau, \theta) \frac{d \varepsilon d \xi}{4 \pi^{2}}$.

We note that (6.56) is asymmetric with respect of the change $\varepsilon$ to $\xi$, however this asymmetry disappears after taking the real part of the expression. Due to this, there is no need for special consideration of the diagram symmetric with respect to the interchange of $\langle\Psi(t)|$ and $|\Psi(t)\rangle$ to the one shown in Fig. 6.4(b).

## The Case of Inverted Ordering

We now turn to the case $\left\langle V_{k m} V_{k m}\right\rangle=\left(V^{\prime}\right)^{2} \neq 0$, and consider the general case $V \neq V^{\prime}$, assuming that a real $V^{\prime}$ is different from another real number $V$. The main contribution comes from the diagrams of the type shown in Fig. 6.4(b), and reads

$$
\begin{align*}
x(t)= & \langle\Psi(t)| \widehat{x}|\Psi(t)\rangle=\operatorname{Re} \sum_{n, m} \int e^{i(\varepsilon-\xi) t}\left\langle x_{m n} V_{m n}\right\rangle \\
& \times \widehat{C}_{m n}(\varepsilon, \xi, \tau, \theta) \widehat{X}_{n}^{-1}(\varepsilon, \tau) \widehat{\rho}_{n}(\varepsilon, \xi, \tau, \theta) \frac{d \varepsilon d \xi}{4 \pi^{2}}, \tag{6.57}
\end{align*}
$$

where $\widehat{C}_{m n}$ denotes an operator similar to $\widehat{\rho}_{n}$ in (6.37) which corresponds to the series Fig. 6.5, and the order of operations is in accordance with that shown in Fig. 6.4(b). The operator $\widehat{X}_{n}^{-1}(\varepsilon, \tau)$ compensates for the repetition of the operator $\widehat{X}_{n}(\varepsilon, \tau)$ found both at the end of $\widehat{\rho}_{n}$ and the beginning of $\widehat{C}_{m n}$.


Fig. 6.5.

We note that the perturbation operators originating from the perturbation series for $\langle\Psi(t)|$ are ordered in a sequence opposite to those of $|\Psi(t)\rangle$ in the "Cooperon" loop $\widehat{C}$ of the diagram in Fig.6.4(b). This may be accounted for when performing calculations for $\widehat{C}$ by analogy to (6.37) for $\widehat{\rho}$ without rewriting all of the cumbersome intermediate expressions. All that is necessary is to change the sign of the derivative over $\theta$ in all of the operators related to $\widehat{C}$. Those operators with inverted ordering are indicated by primes and by analogy to (6.35)-(6.37) we can write

$$
\begin{align*}
\widehat{X}_{n}^{\prime}(\xi, \theta)= & \frac{1}{\xi-\Delta_{n}+i \frac{\partial}{\partial \theta}-i \pi g\left\langle V^{2}\right\rangle q^{2}(\theta)},  \tag{6.58}\\
\widehat{\Xi}^{\prime}(\varepsilon, \xi, \tau, \theta)= & \frac{2 i \pi g\left(V^{\prime}\right)^{2}}{\varepsilon-\xi-i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left(q^{2}(\tau)+q^{2}(\theta)\right)} q(\tau) q(\theta),  \tag{6.59}\\
\widehat{C}_{m n}(\varepsilon, \xi, \tau, \theta)= & \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}^{\prime}(\xi, \theta) \delta_{n}^{m}+\widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}^{\prime}(\xi, \theta)\left(V^{\prime}\right)^{2} \\
& \times q(\tau) q(\theta) \frac{1}{1-\widehat{\widehat{\Xi}}^{\prime}(\varepsilon, \xi)} \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}^{\prime}(\xi, \theta) . \tag{6.60}
\end{align*}
$$

The reasoning behind changing the sign of the derivative becomes evident when we recall the physical image associated with the technique discussed
on p. 245. The derivatives $\frac{\partial}{\partial \theta}$ and $\frac{\partial}{\partial \tau}$ appear when we replace the timedependent system by the time-independent system with a dimensionality higher by unity, and assume that the system moves along this new degree of freedom with a constant velocity, while the Hamiltonian depends on this new coordinate. The opposite ordering therefore implies motion along the same coordinate in the opposite direction.

Equation (6.57) takes the form

$$
\begin{equation*}
x(t)=\operatorname{Re}\langle x V\rangle \sum_{n, m} \int e^{i(\varepsilon-\xi) t} \widehat{C}_{m n}(\varepsilon, \xi, \tau, \theta) \widehat{X}_{n}^{-1}(\varepsilon, \tau) \widehat{\rho}_{n}(\varepsilon, \xi, \tau, \theta) \frac{d \varepsilon d \xi}{4 \pi^{2}}, \tag{6.61}
\end{equation*}
$$

where, again, the operator $\widehat{X}_{n}^{-1}(\varepsilon, \tau)$ is needed to avoid double counting of the operator $\widehat{X}_{n}(\varepsilon, \tau)$ originating from the beginning of the series for $\widehat{C}_{m n}$ and that coming from the termination of the diagrams for $\widehat{\rho}_{n}$.

## Factorization of the Diagrams

Equations (6.58)-(6.61) still remain rather cumbersome operator expressions, ones not very convenient for calculations. In order to simplify, one can employ the operator identity

$$
\begin{align*}
\widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}(\xi, \theta)= & {\left[\widehat{X}_{m}^{-1}(\xi, \theta)-\widehat{X}_{m}^{-1}(\varepsilon, \tau)\right]^{-1}\left[\widehat{X}_{m}(\varepsilon, \tau)-\widehat{X}_{m}(\xi, \theta)\right] } \\
= & \frac{-1}{\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left(q^{2}(\tau)+q^{2}(\theta)\right)} \\
& \times\left[\widehat{X}_{m}(\varepsilon, \tau)-\widehat{X}_{m}(\xi, \theta)\right] \tag{6.62}
\end{align*}
$$

which is easy to prove with the help of (6.41) and the evident relation

$$
\begin{equation*}
\frac{1}{\widehat{X}_{m}(\xi, \theta)}-\frac{1}{\widehat{X}_{m}(\varepsilon, \tau)}=\frac{\widehat{X}_{m}(\varepsilon, \tau)-\widehat{X}_{m}(\xi, \theta)}{\widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}(\xi, \theta)} \tag{6.63}
\end{equation*}
$$

and which allows one to replace the products $\widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}(\xi, \theta)$ of the mean Greens functions by their differences $\widehat{X}_{m}(\varepsilon, \tau)-\widehat{X}_{m}(\xi, \theta)$ subject to the action of an $m$-independent operator. From (6.36) one also sees that

$$
\begin{equation*}
\frac{2 i \pi g}{\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left(q^{2}(\tau)+q^{2}(\theta)\right)}=\sum_{n} \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \tag{6.64}
\end{equation*}
$$

and hence (6.62) reads

$$
\begin{equation*}
\widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}(\xi, \theta)=\frac{-i}{2 \pi g} \sum_{n} \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta)\left[\widehat{X}_{m}(\varepsilon, \tau)-\widehat{X}_{m}(\xi, \theta)\right] \tag{6.65}
\end{equation*}
$$

Equation (6.65) allows one to simplify the triple products of the mean Greens functions to products of independent binary products, such as

$$
\begin{align*}
\sum_{m} \widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}(\xi, \theta) \widehat{X}_{m}^{\prime}(\xi, \theta)= & \frac{-i}{2 \pi g} \sum_{n} \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \\
& \times \sum_{m}\left[\widehat{X}_{m}(\varepsilon, \tau)-\widehat{X}_{m}(\xi, \theta)\right] \widehat{X}_{m}^{\prime}(\xi, \theta)  \tag{6.66}\\
= & \frac{-i}{2 \pi g} \sum_{n} \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \sum_{m} \widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}^{\prime}(\xi, \theta)
\end{align*}
$$

for instance, which gives a very convenient tool for diagram factorization. In the last equation we have taken into account the fact that

$$
\begin{equation*}
\sum_{m} \widehat{X}_{m}(\xi, \theta) \widehat{X}_{m}^{\prime}(\xi, \theta)=0 \tag{6.67}
\end{equation*}
$$

since both $\widehat{X}_{m}(\xi, \theta)$ and $\widehat{X}_{m}^{\prime}(\xi, \theta)$, as functions of $\Delta_{m}$, have poles in the upper part of the complex plane. In Fig. 6.6 we illustrate the topological meaning of this transformation.


Fig. 6.6.

Application of (6.66) to (6.61) yields

$$
\begin{equation*}
x(t)=\operatorname{Re} \frac{-i\langle x V\rangle}{2 \pi g} \int e^{i(\varepsilon-\xi) t} \sum_{n^{\prime}, m} \widehat{C}_{m n^{\prime}}(\varepsilon, \xi, \tau, \theta) \sum_{n} \hat{\rho}_{n}(\varepsilon, \xi, \tau, \theta) \frac{d \varepsilon d \xi}{4 \pi^{2}}, \tag{6.68}
\end{equation*}
$$

which means that the operators $\widehat{C}$ and $\widehat{\rho}$ do not retain common indices. We note that the interaction operator $V$ coupled to the operator $x$ by the ensemble average $\langle x V\rangle$ has been factored out from the integral. This means that we consider the time-dependent susceptibility of the system to the perturbation $V$, which yields the response only after being multiplied by the time dependence of this perturbation. We will occasionally refer to this value as a "response", since it indeed gives a linear response to a small $\delta$-like additional increase in the interaction.

## The Total Population Operator

Considering (6.37) which includes the first term $\widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}(\xi, \theta) \delta_{n}^{0}$ that has been neglected during consideration of the population transfer, and with the allowance of (6.45) and (6.64) one finds

$$
\begin{align*}
& \sum_{n} \rho_{n}(\varepsilon, \xi, \tau, \theta)=\widehat{Z}(\xi, \theta, \varepsilon, \tau)  \tag{6.69}\\
& {\left[\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left(q^{2}(\tau)+q^{2}(\theta)\right)\right] \widehat{X}_{0}(\varepsilon, \tau) \widehat{X}_{0}(\xi, \theta),}
\end{align*}
$$

which can be reduced to a simpler form

$$
\begin{equation*}
\sum_{n} \rho_{n}(\varepsilon, \xi, \tau, \theta)=\widehat{Z}(\xi, \theta, \varepsilon, \tau)\left[\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)\right], \tag{6.70}
\end{equation*}
$$

with the help of the expression

$$
\begin{align*}
{\left[\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\right.} & \left.\left(q^{2}(\tau)+q^{2}(\theta)\right)\right]  \tag{6.71}\\
& \times \widehat{X}_{0}(\varepsilon, \tau) \widehat{X}_{0}(\xi, \theta)=\left[\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)\right],
\end{align*}
$$

which is a particular form of (6.62) for $m=0$. Note that since the operator $\widehat{Z}$ depends only on the difference $\zeta=\varepsilon-\xi$, and hence (6.70), integrated over the variable $\eta=(\varepsilon+\xi) / 2$, has the form

$$
\begin{equation*}
\sum_{n} \rho_{n}(\zeta, \tau, \theta)=\widehat{Z}(\xi, \theta, \varepsilon, \tau) 2 \pi i \tag{6.72}
\end{equation*}
$$

which we will employ later on.
Substitution of the explicit expressions (6.48) into (6.70) yields

$$
\begin{align*}
& \sum_{n} \rho_{n}(\varepsilon, \xi, \tau, \theta)=\delta\left(\frac{\tau-\theta-\tau^{\prime}+\theta^{\prime}}{2}\right) e^{-\pi g V^{2} \int_{\tau^{\prime \prime}}^{\tau} q^{2}(x) d x-\pi g V^{2} \int_{\theta^{\prime \prime}}^{\theta} q^{2}(x) d x} \\
& i \varepsilon\left(\tau-\tau^{\prime \prime}\right)-i \xi\left(\theta-\theta^{\prime \prime}\right)+2 \pi g V^{2} \int_{\frac{\tau+\theta}{2}}^{\frac{\tau+\theta}{2}}\left(q\left(\frac{\tau-\theta}{2}+x\right) q\left(x-\frac{\tau-\theta}{2}\right)\right) d x \\
& e^{\frac{\tau^{\prime}}{2}} \Theta\left(\tau-\tau^{\prime}\right) \Theta\left(\theta-\theta^{\prime}\right)  \tag{6.73}\\
& {\left[\delta\left(\tau^{\prime}-\tau^{\prime \prime}\right) \Theta\left(\theta^{\prime}-\theta^{\prime \prime}\right)+\delta\left(\theta^{\prime}-\theta^{\prime \prime}\right) \Theta\left(\tau^{\prime}-\tau^{\prime \prime}\right)\right] .}
\end{align*}
$$

Note that the Fourier transform of the operator has for $\tau=\theta=t$, a $\delta$-like kernel

$$
\begin{equation*}
\left.\int e^{i(\varepsilon-\xi) t} \sum_{n} \widehat{\rho}_{n}(\varepsilon, \xi, \tau, \theta) \frac{d \varepsilon d \xi}{4 \pi^{2}}\right|_{\tau=\theta=t} \rightarrow \delta\left(\theta^{\prime \prime}\right) \delta\left(\tau^{\prime \prime}\right) \tag{6.74}
\end{equation*}
$$

which is consistent with the normalization requirement

$$
\begin{equation*}
\left.\int_{-\infty}^{\infty} d \tau^{\prime \prime} d \theta^{\prime \prime} \int e^{i(\varepsilon-\xi) t} \sum_{n} \widehat{\rho}_{n}(\varepsilon, \xi, \tau, \theta) \frac{d \varepsilon d \xi}{4 \pi^{2}}\right|_{\tau=\theta=t} \rho\left(t=0, \tau^{\prime \prime}, \theta^{\prime \prime}\right)=1 \tag{6.75}
\end{equation*}
$$

following the application of the operator (6.74) to the $\tau$-independent initial condition $\rho(t=0, \tau, \theta)=1$.

## Sum of the "Cooperon" Series

Substitution of the combination

$$
\begin{equation*}
\widehat{\Xi}^{\prime}(\varepsilon, \xi, \tau, \theta)=\sum_{n} \widehat{X}_{n}(\varepsilon, \tau) \widehat{X}_{n}^{\prime}(\xi, \theta)\left(V^{\prime}\right)^{2} q(\tau) q(\theta) \tag{6.76}
\end{equation*}
$$

entering (6.59) into (6.60) yields, after taking the sum over $n^{\prime}$ and $m$

$$
\begin{align*}
& \sum_{n^{\prime}, m} \widehat{C}_{n^{\prime} m}(\varepsilon, \xi, \tau, \theta)=\sum_{n^{\prime}, m} \widehat{X}_{n^{\prime}}(\varepsilon, \tau) \widehat{X}_{n^{\prime}}^{\prime}(\xi, \theta) \delta_{n}^{m}  \tag{6.77}\\
& +\sum_{n^{\prime}} \widehat{X}_{n^{\prime}}(\varepsilon, \tau) \widehat{X}_{n^{\prime}}^{\prime}(\xi, \theta)\left(V^{\prime}\right)^{2} q(\tau) q(\theta) \frac{1}{1-\widehat{\Xi}^{\prime}} \sum_{m} \widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}^{\prime}(\xi, \theta) \\
= & \left(1+\frac{\widehat{\Xi}^{\prime}}{1-\widehat{\widehat{E}}^{\prime}}\right) \sum_{m} \widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}^{\prime}(\xi, \theta)=\frac{1}{1-\widehat{\widehat{E}}^{\prime}} \sum_{m} \widehat{X}_{m}(\varepsilon, \tau) \widehat{X}_{m}^{\prime}(\xi, \theta) .
\end{align*}
$$

With the allowance of (6.64) and (6.59) this results in

$$
\begin{align*}
\widehat{Z}^{\prime}(\xi, \theta, \varepsilon, \tau) & =\sum_{n^{\prime}, m} \widehat{C}_{n^{\prime} m}(\varepsilon, \xi, \tau, \theta) \\
& =\frac{2 i \pi g}{\varepsilon-\xi-i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left((q(\tau)-q(\theta))^{2}+\lambda q(\tau) q(\theta)\right)}, \tag{6.78}
\end{align*}
$$

where $\lambda=2\left(V^{2}-\left(V^{\prime}\right)^{2}\right) / V^{2}$ is the relative difference between the mean squared matrix elements of the interaction and their mean moduli squared. For real matrix elements, $\lambda$ vanishes.

By analogy to (6.47) one finds the kernel $Z^{\prime}\left(\xi, \theta, \theta^{\prime}, \varepsilon, \tau, \tau^{\prime}\right)$ of this integral operator by solving the differential equation

$$
\begin{align*}
& {\left[\varepsilon-\xi+i \pi g V^{2}\left((q(\tau)-q(\theta))^{2}+\lambda q(\tau) q(\theta)\right)\right] Z^{\prime}\left(\xi, \theta, \theta^{\prime}, \varepsilon, \tau, \tau^{\prime}\right) } \\
+ & i\left(\frac{\partial}{\partial \tau}-\frac{\partial}{\partial \theta}\right) Z^{\prime}\left(\xi, \theta, \theta^{\prime}, \varepsilon, \tau, \tau^{\prime}\right)=2 i \pi g \delta\left(\tau-\tau^{\prime}\right) \delta\left(\theta-\theta^{\prime}\right) \tag{6.79}
\end{align*}
$$

which yields

$$
\begin{gather*}
Z^{\prime}\left(\xi, \theta, \theta^{\prime}, \varepsilon, \tau, \tau^{\prime}\right)=2 \pi g \Theta\left(\tau-\theta-\tau^{\prime}+\theta^{\prime}\right) \delta\left(\frac{\tau+\theta-\tau^{\prime}-\theta^{\prime}}{2}\right) \\
e^{i(\varepsilon-\xi) \frac{\tau-\theta-\tau^{\prime}+\theta^{\prime}}{2}-\pi g V^{2}} \int_{\frac{\tau-\theta}{2}-\theta^{\prime}}^{\frac{\tau-\theta}{2}}\left[\left(q\left(\frac{\tau+\theta}{2}+x\right)-q\left(\frac{\tau+\theta}{2}-x\right)\right)^{2}+\lambda q\left(\frac{\tau+\theta}{2}+x\right) q\left(\frac{\tau+\theta}{2}-x\right)\right] d x \tag{6.80}
\end{gather*}
$$

## The Response

We are now in a position to determine the response. To this end we substitute (6.70), (6.78) into (6.68) and obtain

$$
\begin{align*}
x(t)= & \operatorname{Re} \frac{-i\langle x V\rangle}{2 \pi g} \int e^{i(\varepsilon-\xi) t} \widehat{Z}^{\prime}(\xi, \theta, \varepsilon, \tau) \\
& \widehat{Z}(\xi, \theta, \varepsilon, \tau)\left[\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)\right] \frac{d \varepsilon d \xi}{4 \pi^{2}} . \tag{6.81}
\end{align*}
$$

We note that both $\widehat{Z}^{\prime}(\xi, \theta, \varepsilon, \tau)$ and $\widehat{Z}(\xi, \theta, \varepsilon, \tau)$ depend only on the variable $\zeta$ given by the difference $\varepsilon-\xi=\zeta$, whereas the functional dependence on the sum $\varepsilon+\xi=2 \eta$ enters only via the combination $\widehat{S}=\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)$, related to the initial condition. This immediately results in

$$
\begin{equation*}
x(t)=\operatorname{Re} \frac{-i\langle x V\rangle}{2 \pi g} \int e^{i(\varepsilon-\xi) t} \widehat{Z}^{\prime}(\zeta, \theta, \tau) \widehat{Z}(\zeta, \theta, \tau) \widehat{S}(\zeta, \eta) \frac{d \eta d \zeta}{4 \pi^{2}} . \tag{6.82}
\end{equation*}
$$

From (6.35) one finds

$$
\begin{align*}
\int_{-\infty}^{\infty} \widehat{S}(\zeta, \eta) d \eta & =\int_{-\infty}^{\infty} \frac{d \eta}{\eta-\frac{c}{2}-i \frac{\partial}{\partial \theta}-i \pi g\left\langle V^{2}\right\rangle q^{2}(\theta)} \\
& -\int_{-\infty}^{\infty} \frac{d \eta}{\eta-\frac{\varsigma}{2}+i \frac{\partial}{\partial \tau}+i \pi g\left\langle V^{2}\right\rangle q^{2}(\tau)}=2 \pi i \tag{6.83}
\end{align*}
$$

which is easy to obtain by displacement of the integration contour in the lower part of the complex plane toward $-i \infty$ for the first integral, and in the upper part of the complex plane towards $i \infty$ for the second integral, where the corresponding integrands have no singularities because of the causality principle. The integrals therefore amount to the accumulated phases $\pm i \pi$ of the corresponding logarithms. This yields

$$
\begin{align*}
x(t)= & \operatorname{Re} \frac{\langle x V\rangle}{4 \pi^{2} g} \int e^{i(\varepsilon-\xi) t} \widehat{Z}^{\prime}(\zeta, \theta, \tau) \widehat{Z}(\zeta, \theta, \tau) d \zeta \\
= & \left.\operatorname{Re} \int Z^{\prime}\left(\xi, \theta^{\prime \prime}, \theta, \varepsilon, \tau, \tau^{\prime \prime}\right) Z\left(\xi, \theta^{\prime \prime}, \theta^{\prime}, \varepsilon, \tau^{\prime \prime}, \tau^{\prime}\right)\right|_{\tau=\theta=t} \\
& e^{i(\varepsilon-\xi) t} \frac{\langle x V\rangle}{4 \pi^{2} g} d \zeta d \tau^{\prime \prime} d \theta^{\prime \prime} d \tau^{\prime} d \theta^{\prime} . \tag{6.84}
\end{align*}
$$

where we have taken into account that the convolution of the kernels of $\widehat{Z}^{\prime}$ and $\widehat{Z}$ has to be performed with an inverted order of $\theta$-arguments in the operator $\widehat{Z}^{\prime}$, as illustrated in Fig. 6.7.

We now make use of (6.48) and (6.80) for the kernels of $\widehat{Z}^{\prime}$ and $\widehat{Z}$ which can be written in the form

$$
\begin{align*}
Z^{\prime}\left(\xi, \theta, \theta^{\prime}, \varepsilon, \tau, \tau^{\prime}\right)= & -i \delta\left(\frac{\tau+\theta-\tau^{\prime}-\theta^{\prime}}{2}\right) \Theta\left(\tau-\tau^{\prime}\right) \\
& e^{i \zeta\left(\tau-\tau^{\prime}\right)-\pi g V^{2} \int_{\tau^{\prime}}^{\tau}\left((q(x)-q(\tau+\theta-x))^{2}+\lambda q(x) q(\tau+\theta-x)\right] d x}, \\
Z\left(\xi, \theta, \theta^{\prime}, \varepsilon, \tau, \tau^{\prime}\right)= & 2 \pi g \delta\left(\frac{\tau-\theta-\tau^{\prime}+\theta^{\prime}}{2}\right) \Theta\left(\tau-\tau^{\prime}\right) \\
& e^{i \zeta\left(\tau-\tau^{\prime}\right)-\pi g V^{2} \int_{\tau^{\prime}}^{\tau}(q(x)-q(x-\tau+\theta))^{2} d x} \tag{6.85}
\end{align*}
$$

and arrive at


Fig. 6.7. The correspondence of the arguments $\tau$ and $\theta$ of the operators $\hat{Z}$ and $\hat{Z}^{\prime}$ in (6.84) follows from the topology of the diagrams shown in Fig.6.6.

$$
\begin{align*}
x(t)= & \operatorname{Re} \frac{\langle x V\rangle}{2 \pi i} \int d \zeta d \tau^{\prime \prime} d \theta^{\prime \prime} d \tau^{\prime} d \theta^{\prime} e^{-i \zeta t} \Theta\left(\tau-\tau^{\prime \prime}\right) \delta\left(\frac{\tau+\theta^{\prime \prime}-\tau^{\prime \prime}-\theta}{2}\right) \\
& e^{i \zeta\left(\tau-\tau^{\prime \prime}\right)-\pi g V^{2} \int_{\tau^{\prime \prime}}^{\tau}\left[\left(q(x)-q\left(\tau+\theta^{\prime \prime}-x\right)\right)^{2}+\lambda q(x) q\left(\tau+\theta^{\prime \prime}-x\right)\right] d x} \Theta\left(\tau^{\prime \prime}-\tau^{\prime}\right) \\
& \left.\delta\left(\frac{\tau^{\prime \prime}-\theta^{\prime \prime}-\tau^{\prime}+\theta^{\prime}}{2}\right) e^{i \zeta\left(\tau^{\prime \prime}-\tau^{\prime}\right)-\pi g V^{2} \int_{\tau^{\prime}}^{\tau^{\prime \prime}}\left(q(x)-q\left(x-\tau^{\prime}+\theta^{\prime}\right)\right)^{2} d x}\right|_{\tau=\theta=t} \tag{6.86}
\end{align*}
$$

Integration over $d \zeta$ yields one more $\delta$-function $2 \pi \delta\left(\tau-\tau^{\prime \prime}-t\right)$, and after straightforward integration over $d \theta^{\prime \prime} d \tau^{\prime} d \theta^{\prime}$ of the expression with three $\delta$ functions one finds

$$
\begin{align*}
x(t)= & \operatorname{Re} 2 i\langle x V\rangle \int_{0}^{t} d \tau^{\prime \prime} e^{-\pi g \lambda V^{2} \int_{0}^{\tau^{\prime \prime}} q^{2}(x) d x} \\
& \exp \left\{-\pi g V^{2} \int_{\tau^{\prime \prime}}^{t}\left(q(x)-q\left(t+\tau^{\prime \prime}-x\right)\right)^{2} d x\right\}, \tag{6.87}
\end{align*}
$$

where we have substituted $\tau=\theta=t$, as required in this approach.
Note that for an interaction profile $q(t)$ ranging not from $t=0$ to $t=\infty$, but from $t=-\infty$ to $t=\infty$, an evident generalization of this expression reads

$$
\begin{align*}
x(t)= & \operatorname{Re} 2 i\langle x V\rangle \int_{-\infty}^{t} d \tau^{\prime \prime} e^{-\pi g \lambda V^{2}} \int_{-\infty}^{\tau^{\prime \prime}} q^{2}(x) d x \\
& \exp \left\{-\pi g V^{2} \int_{\tau^{\prime \prime}}^{t}\left(q(x)-q\left(t+\tau^{\prime \prime}-x\right)\right)^{2} d x\right\} . \tag{6.88}
\end{align*}
$$

This short and elegant expression (6.88) has only one disadvantage - it vanishes for a generic quantum system, since all of the averages $V^{2},\langle x V\rangle$, and
$\lambda$ are normally real. Therefore a generic complex quantum system with a uniform spectrum does not manifest either the linear response nor any of the higher responses of odd orders.


Fig. 6.8.

The situation can be different for the even-order responses, if the triple means $\left\langle x_{n m} V_{m k} V_{k n}\right\rangle=\left\langle x V^{2}\right\rangle$ differ from zero. This is, for example, the case for the interaction $\widehat{V}=\mathcal{E}_{x} \widehat{d}_{x}+\mathcal{B}_{y} \widehat{\mu}_{y}$, proportional to the electric field $\mathcal{E}$ in the $x$-direction and the magnetic field $\mathcal{B}$ in the $y$-direction and the operator of the Hall current $\widehat{x}=\widehat{j}_{z}$ in the $z$-direction. In this case the average $\langle x V\rangle$ entering the perturbation series shown in Fig. 6.6 has to be replaced by the mean value

$$
\begin{align*}
\langle x V\rangle & \rightarrow \sum_{k}\left\langle x_{n m} V_{m k} \widehat{X}_{k}^{\prime}(\xi, \theta) V_{k n}\right\rangle \\
& =\left\langle x V^{2}\right\rangle \sum_{k} \widehat{X}_{k}^{\prime}(\xi, \theta)=\left\langle x V^{2}\right\rangle i \pi g \tag{6.89}
\end{align*}
$$

where we have employed (4.32) for $Q(\xi)=\sum_{k} \widehat{X}_{k}(\xi)$. We illustrate the meaning of this replacement and subsequent factorization of the series in Fig. 6.8. This yields

$$
\begin{align*}
x(t)= & -\pi g\left\langle x V^{2}\right\rangle \int_{-\infty}^{t} d \tau^{\prime \prime} e^{-\pi g \lambda V^{2} \int_{-\infty}^{\tau^{\prime \prime}} q^{2}(x) d x} \\
& \exp \left\{-\pi g V^{2} \int_{\tau^{\prime \prime}}^{t}\left(q(x)-q\left(t+\tau^{\prime \prime}-x\right)\right)^{2} d x\right\}, \tag{6.90}
\end{align*}
$$

for the response, which in the dimensionless units $\pi g V^{2} \rightarrow 1, \pi g\left\langle x V^{2}\right\rangle \rightarrow \mathcal{X}$ takes the simple form

$$
\begin{equation*}
x(t)=-\mathcal{X} \int_{-\infty}^{t} d \tau^{\prime \prime} e^{-\lambda \int_{-\infty}^{\tau^{\prime \prime}} q^{2}(x) d x-\int_{\tau^{\prime \prime}}^{t}\left(q(x)-q\left(t+\tau^{\prime \prime}-x\right)\right)^{2} d x} \tag{6.91}
\end{equation*}
$$

It is worth mentioning that, for a complex system without central symmetry, the even order response given by the average $\sum_{k}\left\langle V_{n m} x_{m k} \widehat{X}_{k}^{\prime}(\xi, \theta) V_{k n}\right\rangle$ can also contribute to the result.

In order to gain a first insight into the physics of the response of complex quantum systems, let us consider the simplest step-like time dependence $q(t)=\Theta(t)$, which allows one to write (6.91) in the form

$$
\begin{equation*}
x(t)=-\mathcal{X} \int_{0}^{t} d \tau^{\prime \prime} e^{-\lambda \int_{0}^{\tau^{\prime \prime}} q^{2}(x) d x}=\mathcal{X} \frac{1-e^{-\lambda t}}{\lambda} \tag{6.92}
\end{equation*}
$$

This shows that the response approaches its stationary value $\mathcal{X} / \lambda$ with the rate $\lambda$ proportional to the difference $V^{2}-\left(V^{\prime}\right)^{2}$ given by the mean square $(\operatorname{Im} V)^{2}$ of the imaginary part of the coupling matrix elements. Note that the stationary value itself is inversely proportional to the imaginary part of the coupling, since $\mathcal{X} \sim\left\langle x V^{2}\right\rangle \sim\langle x \operatorname{Re} V \operatorname{Im} V\rangle$.

### 6.3 Harmonic Perturbation of Complex Systems

One of the standard characteristics of a physical system is its behavior under harmonic perturbation. Consideration of Sect.6.2 enables us to give an immediate answer to the question concerning the population distribution and the response of a system with a uniform density of states subjected to a harmonic perturbation proportional to a random matrix. Analysis of this case will be followed by the consideration of several more sophisticated examples of quantum systems with interaction matrices of a block structure. The treatment of these cases is facilitated by the use of the quasi-energy technique (Sect.5.1), to be presented below in more detail.

### 6.3.1 Population Distribution over a Uniform Spectrum

Equations (6.49) and (6.55) allow one to explicitly find the time-dependent population distribution and the response of complex quantum systems for the time-dependent interaction $\chi(t) \cos (\omega t+\varphi) \widehat{V}$, proportional to a random matrix $\widehat{V}$. To shorten the expressions that follow we set $\omega=1,2 \pi g V^{2}=W$ and first consider the population dynamics.

## Population Dynamics for an Abrupt Switching

In the case $\varphi=0$ for the step-like envelope $\chi(t)=\Theta(t)$ the population density reads

$$
\begin{equation*}
\left.\pi g \rho_{n}(t)=\operatorname{Re} \int_{0}^{t} e^{i \Delta_{n} v-W\left(\int_{0}^{t} \cos ^{2}(x) d x-\right.} \int_{0}^{t-v} \cos (x) \cos (v+x) d x\right) ~ d v \tag{6.93}
\end{equation*}
$$

We perform the integration in the exponent and arrive at

$$
\begin{align*}
\pi g \rho_{n}(t)= & \operatorname{Re} \int_{0}^{t} d y \exp \left\{i \Delta_{n} y-\frac{W}{4}(2 t+\sin 2 t)\right. \\
& \left.+\frac{W}{4}(2(t-y) \cos y+\sin (2 t-y)-\sin y)\right\} \tag{6.94}
\end{align*}
$$

For $W$ small, as compared to the frequency, one can neglect the oscillating terms with constant amplitudes and obtain the integral

$$
\begin{equation*}
g \rho_{n}(t)=\operatorname{Re} \int_{0}^{t} d y \exp \left[i \Delta_{n} y-\frac{W}{2} t+\frac{W}{2}(t-y) \cos y\right], \tag{6.95}
\end{equation*}
$$

which for times exceeding the typical time $1 / W$ of a single transition can be calculated approximately. We first note that for a large number of transitions $W t \gg 1$, the main contribution can come only from the vicinities of the points where $\cos y \simeq 1$, otherwise the contribution is small. This immediately yields $y=2 \pi k+z$ with $z \ll 1$, the integration over $z$ can be extended from $-\infty$ to $\infty$ for all $k$ except of $k=0$, and the population (6.95) takes the form

$$
\begin{align*}
\pi g \rho_{n}(t)= & \operatorname{Re} \int_{0}^{\infty} d z \exp \left[i \Delta_{n} z-\frac{W t}{2} \frac{z^{2}}{2}\right] \\
& +\operatorname{Re} \sum_{k=1}^{\infty} \int_{-\infty}^{\infty} d z \exp \left[i \Delta_{n}(2 \pi k+z)-\frac{W t}{2} \frac{z^{2}}{2}-W \pi k\right] \tag{6.96}
\end{align*}
$$

The integral in the first term is half of the integral from $-\infty$ to $\infty$, and the rest of the terms form a geometric series which results in

$$
\begin{align*}
g \rho_{n}(t) & =\frac{2}{\sqrt{\pi W t}} \exp \left[-\frac{\Delta_{n}^{2}}{W t}\right] \operatorname{Re}\left[\frac{1}{2}+\frac{e^{2 \pi i \Delta_{n}-W \pi}}{1-e^{2 \pi i \Delta_{n}-W \pi}}\right] \\
& =\frac{1}{\sqrt{\pi W t}} \exp \left[-\frac{\Delta_{n}^{2}}{W t}\right] \operatorname{Im} \cot \left(\pi \Delta_{n}+i W / 2\right) \\
& =\frac{1}{\sqrt{\pi W t}} \exp \left[-\frac{\Delta_{n}^{2}}{W t}\right] \sum_{m=-\infty}^{\infty} \frac{W / 2 \pi}{\left(\Delta_{n}-m\right)^{2}+W^{2} / 4} \tag{6.97}
\end{align*}
$$

where we have taken into account the singularities of $\cot \left(\pi \Delta_{n}+i W / 2\right)$ which has simple poles at the points $\Delta_{n}=-i W / 2 \pi+m$, where $m$ is an integer. One sees that the population is located in the $W / 2 \pi$-vicinities of the resonant frequencies $\Delta=m$, whereas the population distribution among these vicinities is a Gaussian of width proportional to the square root of the dimensionless time $\sqrt{W t}$. In Fig. 6.9(a) we illustrate this dependence.


Fig. 6.9. A complex quantum system under harmonic perturbation. (a) Population distribution over the uniform spectrum as a function of time for $W / \omega=0.3$. (b) Population distribution for the adiabatic envelope ( 6.98 ) for $\alpha / \omega^{2}=3 \times 10^{-3}$ (solid line) as compared to the distribution for an abrupt switch (dash line). The typical number of transitions $W t=5$ is taken for both cases. The transition rate $W$ is the same as that in (a).

As seen in Fig. 6.9(a), the population remains localized in the vicinities of the resonances. The width of all these populated vicinities is the same and is equal to the transition probability $W$. Moreover, consideration of the general case $q(t)=\cos (\omega t+\varphi)$ shows that this result does not depend on the initial phase $\varphi$ as long as $W$ is small as compared to the frequency, which demonstrates once again the evident fact that all dependencies with fronts shorter than a typical relaxation time $1 / W$ act as an instantaneous switching. A natural question arises as to whether this width is formed during the development of the population dynamics, or is just a result of an abrupt switching of the interaction at $t=0$.

## Population Dynamics for an Adiabatic Switching

In order to answer this question let us consider the population distribution $\rho_{n}(t=\infty)$ following a quasiperiodic perturbation with a smooth Gaussian envelope $q(t)=e^{-\alpha t^{2}} \cos \omega t$. Equation (6.55) reads

$$
\begin{align*}
\pi g \rho_{n}(\infty)= & \operatorname{Re} \int_{-\infty}^{\infty} e^{i \Delta_{n} v-W} \int_{-\infty}^{-\infty} \exp \left(-2 \alpha x^{2}\right) \cos ^{2}(x) d x \\
& e^{W} \int_{0}^{-\infty} \exp \left(-\alpha x^{2}-\alpha(v+x)^{2}\right) \cos (x) \cos (v+x) d x \tag{6.98}
\end{align*} v,
$$

and for $\omega=1$ yields after integration over $x$

$$
\begin{align*}
\pi g \rho_{n}(\infty)= & \operatorname{Re} \int_{-\infty}^{\infty} d v e^{i \Delta_{n} v} \exp \left[-W \sqrt{\frac{\pi}{8 \alpha}}\left(1-e^{-v^{2} \alpha / 2} \cos v\right)\right] \\
& \exp \left[-e^{-1 / 2 \alpha} \sqrt{\frac{\pi}{8 \alpha}} W\left(1-e^{-v^{2} \alpha / 2}\right)\right] \tag{6.99}
\end{align*}
$$

The last factor in the integrand is exponentially close to unity for $\alpha \rightarrow 0$, and therefore can be neglected, while the exponent of the second factor can be expanded in a power series and gives

$$
\begin{equation*}
\pi g \rho_{n}(\infty)=\operatorname{Re} \int_{-\infty}^{\infty} d v \exp \left[i \Delta_{n} v-W \sqrt{\frac{\pi}{8 \alpha}}\left(1-\cos v+\frac{v^{2} \alpha}{2}\right)\right] \tag{6.100}
\end{equation*}
$$

Here we have also taken into account the fact that $W / \sqrt{\alpha}$ is a large factor and therefore the main contribution arises from the domains where $\cos v \simeq 1$.

The population distribution (6.100) is an inverse Fourier transformed product of two factors $e^{-W v^{2} \sqrt{\pi \alpha / 2} / 4}$ and $e^{-W \sqrt{\pi / 8 \alpha}(1-\cos v)}$ each of which has an explicit inverse Fourier transform - a narrow Gaussian

$$
\begin{equation*}
\int_{-\infty}^{\infty} \frac{d v}{2 \pi} e^{i \Delta_{n} v-v^{2} W \sqrt{\pi \alpha / 32}}=\left(\frac{2}{\alpha \pi^{3} W^{2}}\right)^{1 / 4} e^{-\Delta_{n}^{2} \sqrt{2 / \pi \alpha W^{2}}} \tag{6.101}
\end{equation*}
$$

and a series

$$
\begin{equation*}
\int_{-\infty}^{\infty} \frac{d v}{2 \pi} e^{i \Delta_{n} v-W \sqrt{\pi / 8 \alpha}(1-\cos v)}=\sum_{k} \delta\left(\Delta_{n}-k\right) e^{-W \sqrt{\pi / 8 \alpha}} I_{k}\left(W \sqrt{\frac{\pi}{8 \alpha}}\right) \tag{6.102}
\end{equation*}
$$

with the coefficients given by the modified Bessel functions $I_{k}(x)$, respectively. The last equation can be obtained with the help of the Poisson summation formula (3.89) and the Schlafley integral representation for the Bessel functions.

The result is the convolution of (6.101) and (6.102)

$$
\begin{equation*}
g \rho_{n}(\infty)=\sum_{k}\left(\frac{2 \pi}{\alpha W^{2}}\right)^{1 / 4} e^{-\left(\Delta_{n}-k\right)^{2} \sqrt{2 / \pi \alpha W^{2}}} e^{-\sqrt{\pi W^{2} / 8 \alpha}} I_{k}\left(\sqrt{\frac{\pi W^{2}}{8 \alpha}}\right) \tag{6.103}
\end{equation*}
$$

This has a clear physical meaning. The narrow vicinities of the resonances $\Delta_{n}=k$ of width $\left(\pi \alpha W^{2} / 2\right)^{1 / 4}$ given by the span of the Fourier transformation of the Gaussian envelope are populated according to the weights suggested by the solution $e^{-\sqrt{\pi W^{2} / 8 \alpha}} I_{k}\left(\sqrt{\pi W^{2} / 8 \alpha}\right)$ of the one-dimensional diffusion equation on a grid, while the step of the grid is given by the pumping frequency $(\omega=1)$ as we will see from (8.30) later on. The typical number of transitions $\langle | k\rangle$ amounts to the probability $W$ of a single transition multiplied
by the typical interaction time $\sqrt{\frac{\pi}{2 \alpha}}$. The product $W \sqrt{\frac{\pi}{2 \alpha}}$ plays the same role as the combination $W t$ in the previous case. Implicit in this treatment is the condition $\left(\pi \alpha W^{2} / 2\right)^{1 / 4} g \gg 1$ of sufficient spectral density, which means that the populated vicinity of each resonance contains many of quantum levels of the system. In Fig. 6.9(b) we compare the population distributions for the abrupt and the adiabatic switch of the perturbations.

Now we can give an answer to the question concerning the origin of the population distribution in the vicinities of the resonances. It has a coherent nature since the distribution is given by the Fourier transform of the perturbation envelope, being Lorentzian for an abrupt switching, while narrow and Gaussian for a long adiabatic Gaussian pulse of interaction. This implies that the population distribution near the resonance is formed as a result of a reversible process. Indeed, being a broad $\sim W$ Lorentzian for the middle part of the adiabatic interaction pulse, the distribution shrinks to a narrower by a factor of $\left(W^{2} / \alpha\right)^{1 / 4}$ Gaussian when the pulse is over. We can therefore expect the coherent interference phenomena in the case $\omega \sim W$ when the vicinities of different resonances start to overlap. Later on we will consider such a phenomenon for a two-band quantum system.

### 6.3.2 Response of the Uniform Spectrum to a Harmonic Perturbation

We consider the susceptibility of the quantum system with a uniform spectrum to the periodic perturbation switched either abruptly or adiabatically. We employ (6.90) for a perturbation with frequency $\omega=1$ and $2 \pi g V^{2}=W$, and ignore, as earlier, the time dependence of $\left\langle x V^{2}\right\rangle$.

## Response for an Abrupt Switching

For the abrupt switch one has

$$
\begin{align*}
x(t)= & -\pi g\left\langle x V^{2}\right\rangle \int_{0}^{t} d \tau^{\prime \prime} e^{-\lambda W \int_{0}^{\tau^{\prime \prime}} \cos ^{2}(x) d x} \\
& \exp \left\{-W \int_{\tau^{\prime \prime}}^{t}\left(\cos (x)-\cos \left(t+\tau^{\prime \prime}-x\right)\right)^{2} d x\right\} \tag{6.104}
\end{align*}
$$

for the complex quantum system with a uniform spectrum on a periodic perturbation $\widehat{V} \cos t$ starting at $t=0$. Here we set $\omega=1$. We concentrate here on the case of small $W$ and small $\lambda$, which allows one to retain only the linearly increasing component of the first integral in the exponent and write

$$
\begin{align*}
x(t)= & -\pi g\left\langle x V^{2}\right\rangle \int_{0}^{t} d \tau^{\prime \prime} e^{-\lambda W \tau^{\prime \prime} / 2} \\
& \exp \left\{-2 W \sin ^{2}\left(\frac{t+\tau^{\prime \prime}}{2}\right)\left(t-\tau^{\prime \prime}-\sin \left(t-\tau^{\prime \prime}\right)\right)\right\} \tag{6.105}
\end{align*}
$$

Consider this expression in the limit of long $t$. We note that the last exponent is always less than unity, and therefore the contribution of the domain $\tau^{\prime \prime} \gg 1 / \lambda W$ is negligible, and also we neglect the term $\sin \left(t-\tau^{\prime \prime}\right)$ as compared to the large quantity $t-\tau^{\prime \prime}$ thus arriving at

$$
\begin{equation*}
x(t)=-\pi g\left\langle x V^{2}\right\rangle \int_{0}^{t} d \tau^{\prime \prime} e^{-\lambda W \tau^{\prime \prime} / 2} e^{-2 W\left(t-\tau^{\prime \prime}\right) \sin ^{2}\left[\left(t+\tau^{\prime \prime}\right) / 2\right]}, \tag{6.106}
\end{equation*}
$$

where the main contribution to this integral comes from the domains $t+\tau^{\prime \prime} \simeq$ $2 k \pi$ around zeros of the argument of the second exponent. We therefore set $\tau^{\prime \prime}=2 k \pi-\arg e^{i t}+z$ and split the integration interval into a set of $2 \pi$-long intervals, with the total number of these intervals equal to the integer part of the ratio $t / 2 \pi$, that is $K=\left(t-\arg e^{i t}\right) / 2 \pi$, and arrive at

$$
\begin{align*}
x(t)= & -\pi g\left\langle x V^{2}\right\rangle \sum_{k=0}^{K} \exp [-\lambda W(k \pi)] \\
& \int_{-\pi}^{\pi} d z \exp \left[-2 W\left(t-2 k \pi+\arg e^{i t}-z\right) \sin ^{2}[z / 2]\right] \tag{6.107}
\end{align*}
$$

where the small terms $\lambda W(\arg (\exp (i t))-z)$ have been ignored. We make use of the fact that for the all important $k \leq 1 / W \lambda$, the prefactor in front of $\sin ^{2}[z / 2]$ is large, of the order of $W t$, and therefore we can expand the expression in a Taylor series around $z=0$ and keep only the first nonvanishing quadratic term. We can also set both integration limits and the upper limit of the sum to infinity, that is

$$
\begin{align*}
x(t) & =-\pi g\left\langle x V^{2}\right\rangle \sum_{k=0}^{\infty} e^{-\lambda W k \pi} \int_{-\infty}^{\infty} d z \exp \left[-\frac{W t}{2} z^{2}\right] \\
& =\frac{-\pi g\left\langle x V^{2}\right\rangle}{1-e^{-\lambda W \pi}} \sqrt{\frac{2 \pi}{W t}} \simeq \frac{-g\left\langle x V^{2}\right\rangle}{\lambda W} \sqrt{\frac{2 \pi}{W t}}, \tag{6.108}
\end{align*}
$$

We have thus encountered an interesting universal property of complex systems, namely that the second-order response to an abruptly switched harmonic perturbation decreases as the square root of time. This is remarkable in the fact that the periodic perturbation does not result in the oscillation of the permeability of the system. On the contrary, the spectrum of the permeability is located near zero frequencies and manifests there a non-analytic behavior $\sim 1 / \sqrt{\omega}$, as one can find by performing the Fourier transformation of (6.108). It is also important to note that, as we have mentioned on p. 255 , for calculation of the total response, the multiplication of the $1 / \sqrt{t}$ dependence by the time-dependent part of $\left\langle x V^{2}\right\rangle$ is required.

## Response to an Adiabatic Switching

The long-lasting $1 / \sqrt{t}$ character of the response suggests we consider the case of the adiabatic switching of the harmonic interaction. For the case $q(t)=e^{-\alpha t^{2}} \cos \omega t,(6.90)$ for the response reads

$$
\begin{align*}
x(t)= & -\pi g\left\langle x V^{2}\right\rangle \int_{-\infty}^{t} d \tau^{\prime \prime} e^{-\lambda W} \int_{-\infty}^{\tau^{\prime \prime}} \exp \left(-2 \alpha x^{2}\right) \cos ^{2} x d x \\
& \exp \left\{-W \int_{\tau^{\prime \prime}}^{t}\left(e^{-\alpha x^{2}} \cos x-e^{-\alpha\left(t+\tau^{\prime \prime}-x\right)^{2}} \cos \left(t+\tau^{\prime \prime}-x\right)\right)^{2} d x\right\}, \tag{6.109}
\end{align*}
$$

where the rate $\alpha$ of the front rising is a very small number, and the factor in front of the integral we treat as a constant for the moment. We replace $\cos ^{2} x$ entering the first integral in the exponent by its mean value $1 / 2$ assuming $\lambda W \ll 1$, and note that this integral is large $\sim \lambda W / \sqrt{\alpha}$ unless the integration limit $\tau^{\prime \prime}$ is a large negative number $\tau^{\prime \prime} \ll-1 / \sqrt{\alpha}$. Since the second integral in the exponent is negative, the last inequality is also a necessary condition for the total integral not to vanish. Indeed, for $\tau^{\prime \prime} \sim-1 / \sqrt{\alpha}$ the result of the integration will be small because of the first integral in the exponent, which takes a large negative value. In turn, the second integral in the exponent becomes large when two oscillating terms of the integrand have different amplitudes. This means that the response is exponentially small, unless $t \sim$ $-\tau^{\prime \prime} \gg 1 / \sqrt{\alpha}$. Therefore the sum $t+\tau^{\prime \prime}$ must be a small value, such that $\alpha\left(t+\tau^{\prime \prime}\right)^{2} \ll 1$ and the limits of the second integral in the exponent therefore can be taken as infinite. After straightforward integration of the second term one obtains

$$
\begin{align*}
x(t)= & -\pi g\left\langle x V^{2}\right\rangle \int_{-\infty}^{t} d \tau^{\prime \prime} e^{-\lambda W} \int_{-\infty}^{\tau^{\prime \prime}} \exp \left(-2 \alpha x^{2}\right) / 2 d x \\
& \exp \left\{-W \sqrt{\frac{\pi}{2 \alpha}}\left(1-e^{-\alpha\left(t+\tau^{\prime \prime}\right)^{2} / 2} \cos \left(t+\tau^{\prime \prime}\right)\right)\right\} \tag{6.110}
\end{align*}
$$

where we have also set $e^{-1 / \alpha}$ to zero.
The second exponential factor dominates in (6.110), since it contains a large parameter $W \sqrt{\frac{\pi}{2 \alpha}}$, and the main contributions to the integral arise in the vicinities of the points $\tau^{\prime \prime}=-t+2 k \pi$ periodically arranged on the integration axis, where $\cos \left(t+\tau^{\prime \prime}\right) \rightarrow 1$. Each interval of length $2 \pi$ gives the contribution $\left(\frac{\pi \alpha}{2}\right)^{1 / 4} W^{-1 / 2} e^{\alpha\left(t+\tau^{\prime \prime}\right)^{2} / 4}$, equal to the width of the corresponding vicinity, which one can find by expanding the second exponent in series near these points. Calculation of the integrals (6.110) for small $\alpha$ results in

$$
\begin{align*}
x(t)= & -\frac{\pi g\left\langle x V^{2}\right\rangle}{2 \sqrt{W}}\left(\frac{\pi \alpha}{2}\right)^{1 / 4} \sum_{k=-\infty}^{\infty} e^{-\lambda W} \int_{-\infty}^{-t+2 k \pi} \exp \left(-2 \alpha x^{2}\right) / 2 d x \\
& \exp \left\{\alpha(k \pi)^{2}-W \sqrt{\frac{\pi}{2 \alpha}}\left(1-e^{-2 \alpha(k \pi)^{2}}\right)\right\} . \tag{6.111}
\end{align*}
$$

In order to be consistent within the approximation, one has to neglect the first term in the second exponent as compared to the second one and take the Taylor expansion of $1-e^{-2 \alpha(k \pi)^{2}}$. This yields

$$
\begin{gather*}
x(t)=-\frac{\pi g\left\langle x V^{2}\right\rangle}{2 \sqrt{W}}\left(\frac{\pi \alpha}{2}\right)^{1 / 4} \\
\sum_{k=-\infty}^{\infty} e^{-W \sqrt{2 \alpha \pi}(k \pi)^{2}-\lambda W \int_{-\infty}^{-t+2 k \pi}} \exp \left(-2 \alpha x^{2}\right) / 2 d x \tag{6.112}
\end{gather*}
$$

For small $\alpha$ the sum can be replaced by an integral which results in

$$
\begin{align*}
x(t)= & -\frac{g\left\langle x V^{2}\right\rangle}{4 \sqrt{2 \alpha} \sqrt{W}}\left(\frac{\pi \alpha}{2}\right)^{1 / 4} \int_{-\infty}^{\infty} d y \\
& \exp \left[-\frac{W \sqrt{\pi}}{4 \sqrt{2 \alpha}}\left((y+t \sqrt{2 \alpha})^{2}+\frac{\lambda}{\sqrt{\pi}} \int_{-\infty}^{y} \exp \left(-x^{2}\right) d x\right)\right] \tag{6.113}
\end{align*}
$$

which can be found by the saddle-point method.
For small $\lambda$ the saddle point is located near $y=-t \sqrt{2 \alpha}$, whereas the second derivative is close to $W \sqrt{\pi / 8 \alpha}$ and therefore

$$
\begin{equation*}
x(t)=-\frac{g\left\langle x V^{2}\right\rangle \sqrt{\pi}}{2 \sqrt{2} W} \exp \left[-\frac{W \lambda}{4 \sqrt{2 \alpha}} \int_{-\infty}^{-t \sqrt{2 \alpha}} \exp \left(-x^{2}\right) d x\right] \tag{6.114}
\end{equation*}
$$

which in the notation of (6.91) and after the replacements $\pi g\left\langle x V^{2}\right\rangle \rightarrow \mathcal{X}$ for dimensionless permeability, $W \lambda \sqrt{\pi} / 8 \sqrt{2 \alpha} \rightarrow A$ for the parameter representing a typical number of transitions during the adiabatic pulse, and $t \sqrt{2 \alpha} \rightarrow T$ for dimensionless time takes the short form

$$
\begin{equation*}
x(t)=-\frac{\mathcal{X}}{2 \sqrt{2 \pi}} \exp \{A \operatorname{Erf}(T)-A\} \tag{6.115}
\end{equation*}
$$

where the Gaussian integral is replaced by the error function

$$
\begin{equation*}
\int_{-\infty}^{-X} \exp \left(-x^{2}\right) d x=\frac{\sqrt{\pi}}{2}(1-\operatorname{Erf}[X]) \tag{6.116}
\end{equation*}
$$

that has already been encountered in (3.132).


Fig. 6.10. Response of a complex quantum system (6.117) to a harmonic perturbation switched on adiabatically. (a) The response (arbitrary units) as a function of dimensionless time $T=t \sqrt{2 \alpha}$ and the parameter $A$ proportional to the number of transitions during the pulse $\sim W / \sqrt{\alpha}$ and the ratio $\sqrt{\lambda}=\operatorname{Im} V / \operatorname{Re} V$ of the mean squared imaginary and the mean squared real parts of the interaction. (b) Response as a function of $T$ for $A=9$ (bold line) and for $A=100$ (dashed line). The main part of the response is located near the final edge of the pumping pulse (solid line), when the decay rate drops faster than the pumping. (c) The stronger the interaction the smaller the response, since the decay rate drops to the required level later.

By adding $-T^{2}$ to the exponent one can also take into account the time dependence of $\left\langle x V^{2}\right\rangle$. Discussing the response (6.92), we have already mentioned that $\lambda \sim \operatorname{Im} V^{2}$ and hence $A \sim \lambda$, whereas $\mathcal{X} \sim \operatorname{Im} V \sim \sqrt{A}$. For small $\lambda$, assuming the mean squared real part of the perturbation fixed, we can write

$$
\begin{equation*}
x(T) \sim \sqrt{A} \exp \left\{A \operatorname{Erf}(T)-A-T^{2}\right\} \tag{6.117}
\end{equation*}
$$

for the time-dependent response given by the product of the permeability and the square of the time-dependent Gaussian envelope of the interaction $q(t)$. In Fig. 6.10(a) we illustrate this dependence. The response is mainly located near the end of the pumping pulse, which is a natural consequence of the fact that the decay rate increases rapidly at the pulse front, thus resulting in a rapid decrease of the response, whereas at the end of the pulse the equally rapid decrease of the rate does not last a sufficiently long time to completely destroy the response. We note, that the integral of $\exp \left\{A \operatorname{Erf}(T)-A-T^{2}\right\}$ over $T$ tends for large $A$ to the dependence $0.88 / A$. With dimensional variables, this yields the asymptotic behavior

$$
\begin{equation*}
\int x(t) d t=-\frac{g\left\langle x V^{2}\right\rangle \sqrt{2 \pi} 0.88}{W^{2} \lambda} \sim \frac{x_{t} g}{g W}, \tag{6.118}
\end{equation*}
$$

for the integrated total response, decreasing with the increase of the interaction strength and independent of the width $\alpha^{-1}$ of the pulse envelope.

The numeric factor in (6.118) relates to the chosen model of a Gaussian envelope, but the inverse dependence on the interaction intensity $\sim W^{-1}$ is a general property of the physical quantities given by the self-intersecting diagrams, as is the case in Fig. 6.4 for the response diagrams, since each self-intersection reduces the number of corresponding diagrams by a factor equal to the number of levels in resonance, discussed on p. 75. The order of magnitude estimations suggest that the integrated response is of the order of the typical matrix element $x_{t}$ multiplied by the typical return time $g$ and divided by the number $g W$ of levels in the populated vicinities of the resonances.

### 6.4 Two-Frequency Excitation of Complex Systems

We now turn to complex quantum systems subjected to the simultaneous action of two harmonics and consider the universal properties of nonlinear wave mixing in such systems. We start with the population distribution given by (6.55) for $q(x)=e^{-\alpha t^{2}}(\cos t+\cos b t)$ and then consider the corresponding response (6.87).

### 6.4.1 Population Dynamics for Bi-Harmonic Excitation

After a straightforward calculation one finds the integral

$$
\begin{equation*}
\int_{-\infty}^{\infty}\left[q^{2}(x)-q(x) q(v+x)\right] d x=\frac{1}{2} \sqrt{\frac{\pi}{2 \alpha}}\left[2-e^{-v^{2} \alpha / 2}(\cos v+\cos b v)\right] \tag{6.119}
\end{equation*}
$$

where the terms $\sim \exp [-1 / \alpha]$ have been ignored, and equation (6.55) yields

$$
\begin{equation*}
\pi g \rho_{n}(t)=\operatorname{Re} \int_{-\infty}^{\infty} \exp \left\{i \Delta_{n} v-W \sqrt{\frac{\pi}{8 \alpha}}\left[2-e^{-v^{2} \alpha / 2}(\cos v+\cos b v)\right]\right\} d v \tag{6.120}
\end{equation*}
$$

One recognizes here the Fourier transformation of the product of two similar expressions

$$
\begin{equation*}
\exp \left\{-W \sqrt{\frac{\pi}{8 \alpha}}\left[1-e^{-v^{2} \alpha / 2} \cos v\right]\right\} \tag{6.121}
\end{equation*}
$$

and

$$
\begin{equation*}
\exp \left\{-W \sqrt{\frac{\pi}{8 \alpha}}\left[1-e^{-v^{2} \alpha / 2} \cos b v\right]\right\} \tag{6.122}
\end{equation*}
$$

each of which has exactly the same structure as the expression in (6.99), which we have encountered in considering the harmonic excitation of the complex system. The Fourier transform of a product $F[a b]$ yields the convolution $F[a] * F[b]$ of the Fourier transforms of the factors, and therefore with the help of (6.103)

$$
\begin{align*}
g \rho_{n}(\infty)= & \sum_{k}\left(\frac{2 \pi}{\alpha W^{2}}\right)^{1 / 4} e^{-\left(\Delta_{n}-k\right)^{2} \sqrt{2 / \pi \alpha W^{2}}} e^{-\sqrt{\pi W^{2} / 8 \alpha}} I_{k}\left(\sqrt{\frac{\pi W^{2}}{8 \alpha}}\right) \\
& * \sum_{m}\left(\frac{2 \pi}{\alpha W^{2}}\right)^{1 / 4} e^{-\left(\Delta_{n}-m b\right)^{2} \sqrt{2 / \pi \alpha W^{2}}} e^{-\sqrt{\pi W^{2} / 8 \alpha}} I_{m}\left(\sqrt{\frac{\pi W^{2}}{8 \alpha}}\right) \\
= & \sum_{k m}\left(\frac{2 \pi}{\alpha W^{2}}\right)^{1 / 4} e^{-\left(\Delta_{n}-k-b m\right)^{2} \sqrt{2 / \pi \alpha W^{2}}} e^{-\sqrt{\pi W^{2} / 2 \alpha}} \\
& I_{m}\left(\sqrt{\frac{\pi W^{2}}{8 \alpha}}\right) I_{k}\left(\sqrt{\frac{\pi W^{2}}{8 \alpha}}\right) \tag{6.123}
\end{align*}
$$

Equation (6.123) has a clear physical meaning. It shows that the population is located in narrow ( $\sim \sqrt{W \sqrt{\alpha}}$ ) vicinities of the resonances $\Delta=k+b m$, whereas the total populations $P_{m k}$ of different vicinities are distributed according the solution $e^{-w} I_{m}(w / 2) I_{k}(w / 2)$ of the diffusion equation on a twodimensional grid (see (8.30)), given in terms of the modified Bessel functions and corresponding to the effective number of transitions $w=W \int q^{2}(t) d t=$ $W \sqrt{\pi / 2 \alpha}$. For large $T$ one finds the asymptotic relation

$$
\begin{equation*}
P_{m k}=e^{-w} I_{m}(w / 2) I_{k}(w / 2) \simeq e^{-\left(k^{2}+m^{2}\right) / w} / \pi w \tag{6.124}
\end{equation*}
$$

and (6.123) for the population spectral density $\rho\left(\Delta_{n}\right)=g \rho_{n}(t=\infty)$ adopts the form

$$
\begin{equation*}
\rho(\Delta)=\sum_{k m}\left(\frac{8 \alpha \pi}{W^{2}}\right)^{1 / 4} e^{-(\Delta-k-b m)^{2} \sqrt{2 / \pi \alpha W^{2}}} e^{-\left(k^{2}+m^{2}\right) \sqrt{2 \alpha / \pi W^{2}}} \frac{1}{\pi^{2} W} \tag{6.125}
\end{equation*}
$$

This expression gives the population distribution over the uniform energy spectrum, which results from the overlap of the vicinities of different resonances, located in the same energy positions. The way these vicinities are distributed along the energy axes depends on the ratio $b$ of the pumping frequencies. The situation is similar to the Bragg scattering patterns as a function of different orientations of solids with rectangular lattices. In Fig. 6.11 we illustrate this dependence.

It turns out that the commensurability of two pumping frequencies plays an important role and determines the pattern of the population distribution. Indeed, when the frequency ratio $b=p / q$ is a rational number given by an ir-


Fig. 6.11. Population distribution (6.123) of a complex quantum system subjected to an action of an adiabatically switched two-frequency field. The distribution results from the diffusion on the two-dimensional grid (upper panel) given by the resonant absorption or emission of quanta at frequencies $\omega_{1}=1$ (index $k$ ) and $\omega_{2}=b$ (index $m$ ) of the field. The result depends on the diffusion-weighted $P_{m, k}$ projection of this grid on the one-dimensional energy axis (lower panel). For the commensurate frequencies $b=p / q$ (left panel) and for the non-commensurate frequencies $b \neq p / q$ the distributions are significantly different. The positions of the resonances are shown by arrows.
reducible fraction of two integer numbers $p$ and $q$, then according to (6.123), the energy positions of different resonances $\Delta=k+b m=(k q+p m) / q$ are separated by a minimum distance $1 / q$. The resulting distribution than turns out to be a set of narrow spikes of population separated by unpopulated intervals. The width of the spikes depends only on the spectral resolution $\sim\left(2 / \pi \alpha W^{2}\right)^{-1 / 4}$ corresponding to the transition rate $W$ and the pulse duration $\sim 1 / \sqrt{\alpha}$, whereas the total population of each vicinity depends on the number of steps needed for arriving at this resonance from the initially populated state and the interaction pulse duration.

The situation is completely different for an irrational $b$. For asymptotically long pulse durations this results in a uniform distribution over the energy scale. However, this distribution is not attained at once but arises through a series of quasiperiodical distributions, each of which is given by the denominator of the best rational approximation to the irrational number $b$ corresponding to the actual spectral resolution $\sim\left(2 / \pi \alpha W^{2}\right)^{-1 / 4}$. In Fig. 6.12 we illustrate this situation for the example of $b=1 / 3$ and $b=1 / \pi$. For the last case, the best approximations are $b=1 / \pi \simeq 1 / 3 \simeq 7 / 22 \simeq 113 / 335$ $\simeq 99532 / 312689$, and one sees the periodic structures corresponding to three


Fig. 6.12. Density of population ( see (6.123) and Fig. 6.11) for different pulse durations $1 / \sqrt{\alpha}$ and different frequency ratios $\omega_{1} / \omega_{2}=b$. Transition rate $W=0.1$ for all plots. (a) For commensurate frequencies $b=1 / 3$, the population distribution is given by the spikes located periodically at the position of the inverse of the largest common divider of the oscillation periods, $\Delta_{p}=1 / 3$. The non-commensurate ratio $b=1 / \pi$ yields a multiperiodic structure. (b) The genesis of multiperiodicity for irrational $b$. For different pulse durations the spectral resolution are different, and the period $1 / q$ is given by the denominators of the best rational approximation $b \simeq$ $p / q$ of the frequency ratio. For $\alpha=10^{-4}$ one finds $q=3$ (dotted line), for $\alpha=10^{-7}$ $q=22$ (dashed line scaled $\times 3$ relative to (a)) and $q=355$ for $\alpha=10^{-11}$ (solid line scaled $\times 9$ ), which corresponds the the best rational approximations $\pi \simeq 3 \simeq$ $22 / 7 \simeq 355 / 113$, respectively.
first denominators $1 / q=1 / 3 ; 1 / 22 ; 1 / 355$ that sequentially attains with increasing pulse duration.

### 6.4.2 Response to Bi-Harmonic Excitation

A complex periodic structure, originating in the number theory relating to the best rational approximation of irrational numbers also appears, although in a less pronounced way, in the spectrum for the universal time-dependent response of a complex quantum system with a uniform density of states, in the presence of a two-frequency perturbation. We shall consider the biharmonic time dependence $q(x)=\chi(t)(\cos t+\cos b t)$ with an envelope $\chi(t)$. The particular choice of $\chi(t)$ does not play a crucial role and can be made in a way minimizing technical complications during the calculations by assuming that most of the time the envelope takes a constant value, say unity. The frequency spectrum of the response according to (6.87) reads

$$
\begin{align*}
x(\Delta) \sim & \int_{-\infty}^{\infty} d t \chi^{2}(t) \int_{-\infty}^{0} d y  \tag{6.126}\\
& e^{i \Delta t-\lambda W \int_{-\infty}^{y+t} q^{2}(x) d x-W \int_{y}^{0}(q(x+t)-q(t+y-x))^{2} d x}
\end{align*}
$$

where we have made the replacements $\tau^{\prime \prime} \rightarrow t+y, x \rightarrow x+t$, and by introducing the factor $\chi^{2}(t)$ we have also taken into account the profile of the pumping field. Before taking the Fourier transform over time we note that the integrand can be represented as a product of two factors

$$
\begin{align*}
x(\Delta) & \sim \int_{-\infty}^{\infty} d t e^{i \Delta i} \widetilde{R}(t) \widetilde{f}(t) \\
\widetilde{f}(t) & =\chi^{2}(t) \exp \left[-\lambda W \int_{-\infty}^{t} q^{2}(x) d x\right] \\
\widetilde{R}(t) & =\int_{-\infty}^{0} d y e^{-W \int_{y}^{0}\left[(q(x+t)-q(t+y-x))^{2}+\lambda q^{2}(x+t)\right] d x} \tag{6.127}
\end{align*}
$$

which means that the Fourier transform of the response is a convolution of two profiles: the term

$$
\begin{equation*}
R(\Delta)=\int_{-\infty}^{\infty} \frac{d t}{2 \pi} e^{i \Delta t} \int_{-\infty}^{0} d y e^{-W \int_{y}^{0}\left[(q(x+t)-q(t+y-x))^{2}+\lambda q^{2}(x+t)\right] d x} \tag{6.128}
\end{equation*}
$$

which allows for the frequency dependence of the response on a large scale and which has the structure of a number of $\delta$-like spikes similar to (6.17), and the term

$$
\begin{equation*}
f(\Delta)=\int_{-\infty}^{\infty} \frac{d t}{2 \pi} \chi^{2}(t) e^{i \Delta t-\lambda W} \int_{-\infty}^{t} q^{2}(x) d x \tag{6.129}
\end{equation*}
$$

allowing for the profile (line shape) of each spike that mainly depends on the steepness of fronts and the shape of the envelope $\chi(t)$.

In the limit of small $\lambda$ we can neglect the oscillating component of $q^{2}(x)$ and taking $\chi^{2}(t)=1+\tanh \alpha x$, for instance, obtain after the replacement $t \rightarrow T / \alpha$

$$
\begin{align*}
f(\Delta) & \left.\left.=\int_{-\infty}^{\infty} \frac{d t}{2 \pi}[1+\tanh \alpha t)\right] e^{i \Delta t-\lambda W} \int_{-\infty}^{t}[1+\tanh \alpha x)\right] d x \\
& =\int_{-\infty}^{\infty} \frac{d T}{\pi \alpha} \frac{e^{2 T}}{1+e^{2 T}} \exp \left\{i \frac{\Delta}{\alpha} T-\frac{\lambda W}{\alpha} \ln \left[1+e^{2 T}\right]\right\} \tag{6.130}
\end{align*}
$$

The substitution $z=e^{2 T}$ results in the line shape

$$
\begin{equation*}
f(\Delta)=\int_{0}^{\infty} \frac{z^{i \Delta / 2 \alpha} d z}{2 \pi \alpha(1+z)^{1+\lambda W / \alpha}}=\frac{\Gamma\left(1+i \frac{\Delta}{2 \alpha}\right) \Gamma\left(\frac{\lambda W}{\alpha}-i \frac{\Delta}{2 \alpha}\right)}{2 \pi \lambda W \Gamma\left(\frac{\lambda W}{\alpha}\right)} \tag{6.131}
\end{equation*}
$$

whose absolute value has a bell-like shape of width $\sim \alpha$ centered at $\Delta=0$ and with phase modulations antisymmetric with respect to this point.

Let us turn now to the term (6.128). The integrand in the exponent is always positive, and hence only $y \sim 1 / \lambda W$ can be important. In this domain we can ignore the variation of the envelope and set $\chi(t)=1$. Moreover, we assume $W \ll 1$ and ignore the oscillatory terms of constant amplitude, retaining only the leading terms with amplitudes proportional to $y$. This yields for the integral in the exponent

$$
\begin{align*}
& \int_{y}^{0}\left[(q(x+t)-q(t+y-x))^{2}+\lambda q^{2}(x+t)\right] d x  \tag{6.132}\\
& \simeq-y\{2+\lambda-\cos [2 t+y]-\cos [b(2 t+y)]\}
\end{align*}
$$

and hence

$$
\begin{align*}
R(\Delta) & =\int_{-\infty}^{\infty} \frac{d t}{2 \pi} e^{i \Delta t} \int_{-\infty}^{0} d y e^{W y\{2+\lambda-\cos [2 t+y]-\cos [b(2 t+y)]\}} \\
& =\int_{-\infty}^{0} d y e^{-i \Delta y / 2+W y \lambda} \int_{-\infty}^{\infty} \frac{d t}{2 \pi} e^{i \Delta t} e^{W y\{2-\cos 2 t-\cos 2 b t\}} . \tag{6.133}
\end{align*}
$$

In (6.120) we have already encountered the last integral, which can be given in terms of Bessel functions, and therefore we arrive at

$$
\begin{equation*}
R(\Delta)=\int_{-\infty}^{0} d y e^{W y \lambda-i \Delta y / 2} e^{2 W y} \sum_{m, n} I_{m}(-W y) I_{n}(-W y) \delta(\Delta-n-b m) \tag{6.134}
\end{equation*}
$$

The structure of this expression resembles that for the population distribution (6.125), which is also a set of narrow (here $\delta$-like) spikes weighted with an intensity distribution function $P_{n, m}$. Following convolution with the profile (6.131) this results in

$$
\begin{align*}
R(\Delta) & =\sum_{m, n} P_{n, m}(W, \lambda) f(\Delta-n-b m) \\
P_{n, m} & =\int_{-\infty}^{0} d y e^{W y \lambda-i(n-b m) y / 2} e^{2 W y} I_{m}(-W y) I_{n}(-W y) . \tag{6.135}
\end{align*}
$$

However, in contrast to the population distribution, the weighting function for the response remains essentially different from the Gaussian. When we
make use of the Gaussian asymptotic expression for the Bessel functions (6.124) it takes the form

$$
\begin{align*}
P_{n, m} & =\int_{-\infty}^{0} \frac{2 d y}{\pi W y} \exp \left\{W y \lambda-\frac{i(n+b m) y}{2}+\frac{\left(k^{2}+m^{2}\right) 2}{W y}\right\} \\
& =-\frac{4}{W} K_{0}\left[2 \sqrt{\left(\frac{i(b m-n)}{W}+2 \lambda\right)\left(m^{2}+n^{2}\right)}\right] \tag{6.136}
\end{align*}
$$

where $K_{0}(z)$ is the modified Bessel function.


Fig. 6.13. Frequency dependence of the response (6.137) of a complex quantum system with a uniform density of states to a bi-harmonic excitation. The principal effect occurs around zero frequency. (a) Coalescence of two resonances for a frequency ratio $b$, assuming a rational value. (b) In contrast to the population distribution Fig. 6.11, the contribution of higher-order resonances is suppressed by the spectral profile (dashed line) decreasing exponentially for large detunings. Therefore only a few resonances near zero frequency can interfere as illustrated in the inset. The parameters of the plot are $W=0.1, \lambda=0.05, \alpha=0.01$.

Combination of (6.131) and (6.136) gives the spectral composition of the response

$$
\begin{align*}
x(\Delta)= & 2 \pi g\left\langle x V^{2}\right\rangle \sum_{m, n} K_{0}\left[2 \sqrt{\left(\frac{i(b m-n)}{W}+2 \lambda\right)\left(m^{2}+n^{2}\right)}\right] \\
& \frac{\Gamma\left(1+i \frac{(\Delta-n-b m)}{2 \alpha}\right) \Gamma\left(\frac{\lambda W}{\alpha}-i \frac{(\Delta-n-b m)}{2 \alpha}\right)}{\pi \lambda W^{2} \Gamma\left(\frac{\lambda W}{\alpha}\right)} . \tag{6.137}
\end{align*}
$$

The main difference between this function and the population distribution (6.125) is due to the narrow weighting function (6.136) given by $K_{0}(z)$ exponentially decreasing at large values of the argument $z$. Therefore the typical detuning $(b m-n)$ between the higher-order harmonics $n$ and $m$ should be smaller than $W /\left(m^{2}+n^{2}\right)$, otherwise there is no substantial contribution to the sum (6.137). However a typical detuning scales as $1 / n$ whereas for the best rational approximation it can be estimated as $1 / n^{2}$, assuming a uniform, random distribution of the higher-order resonances on the frequency axis. Therefore for a typical $b$, the minimum value of the combination $(b m-n)\left(m^{2}+n^{2}\right)$ for integer $n$ and $m$ ranging from $-N$ to $N$ does not depend much on the size $N$ of the interval, remaining of the order of unity. For example, in the case $b=1 / \pi$ for the first four best rational approximations this combination takes the values $0.15,-0.068,-0.0037,-0.317$. This means that for a generic ratio of frequencies at certain conditions one may observe the response consisting of several higher order resonances at once, one for each best rational approximation, but it never forms a periodic pattern as was the case for the population distribution. Typically one can see the coalescence of only two resonances when $b$ assumes a rational value, as shown in Fig. 6.13.

### 6.5 Two-Band System in a Periodic Field

The behavior of two-band systems under the action of a harmonic perturbation shown in Fig.6.14(a) is almost identical to the dynamics of the singleband spectra of Sect. 6.3. The main distinction is in the fact that the evenorder resonances correspond to the first band, while the odd-order resonances correspond to the second one. We therefore address a natural question specific for the two-band problem, about the time dependence of the population difference of the two bands. We also consider here the role of the phase statistics of the pumping field and the regime of overlapping resonances two problems that have not been considered earlier, although they could be formulated for a single-band as well. Note that the two band model can be applied to a single band system when we consider in the resonant approximation a high-frequency perturbation $V_{h}=V q(t) \cos \omega_{h} t$ at a frequency $\omega_{h}$ comparable with the total width $\Gamma \sim \omega_{h}$ of the band. In the last case, the oscillation of the field amplitude $V q(t)$ plays the role of the time dependent perturbation coupling two resonant parts, as shown in Fig 6.14(b).

### 6.5.1 Dynamics of Total Band Populations

We address this problem with the help of the technique employed earlier in this chapter, and then propose a somewhat different point of view at the obtained results, which is based on the consideration of a quantized pumping field.

## The Direct Consideration

Equations (4.28), valid for the total population of the bands in the case of time-independent perturbations, can also be employed for the harmonic pumping by simply replacing the static operators $\hat{X}$ (4.21) by similar operators for the time-dependent situation. For two identical bands the simplest forms have the expressions for the sum $\widehat{\rho}_{+}=\widehat{\rho}_{1}+\widehat{\rho}_{2}$ and the difference $\widehat{\rho}_{-}=\widehat{\rho}_{1}-\widehat{\rho}_{2}$ of the populations

$$
\begin{align*}
& \widehat{\rho}_{+}(\varepsilon, \xi, \tau, \theta)=\frac{1}{1-\widehat{\Xi}(\varepsilon, \xi, \tau, \theta)} \widehat{X}_{0}(\varepsilon, \tau) \widehat{X}_{0}(\xi, \theta), \\
& \widehat{\rho}_{-}(\varepsilon, \xi, \tau, \theta)=\frac{1}{1+\widehat{\Xi}(\varepsilon, \xi, \tau, \theta)} \widehat{X}_{0}(\varepsilon, \tau) \widehat{X}_{0}(\xi, \theta), \tag{6.138}
\end{align*}
$$

where $\widehat{X}_{0}(\varepsilon, \tau)$ and $\widehat{\Xi}(\varepsilon, \xi, \tau, \theta)$ are given by (6.35), (6.36).
We note that (6.38) can be generalized in order to include both signs in the denominator

$$
\begin{align*}
\frac{1}{1 \pm \widehat{\Xi}(\varepsilon, \xi)}= & \frac{1}{\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}(q(\tau) \pm q(\theta))^{2}} \\
& {\left[\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}\left(q^{2}(\tau)+q^{2}(\theta)\right)\right] } \tag{6.139}
\end{align*}
$$

which with the allowance of (6.71) results in the equation

$$
\begin{align*}
& {\left[\varepsilon-\xi+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}(q(\tau)+q(\theta))^{2}\right]} \\
& \widehat{\rho}_{-}(\varepsilon, \xi, \tau, \theta)=\left[\widehat{X}_{0}(\xi, \theta)-\widehat{X}_{0}(\varepsilon, \tau)\right] \tag{6.140}
\end{align*}
$$

for the population difference.
After the standard replacements $\varepsilon-\xi \rightarrow \zeta, \varepsilon+\xi \rightarrow 2 \eta$ and integration over $\eta$ this yields

$$
\begin{equation*}
\left[\zeta+i \frac{\partial}{\partial \theta}+i \frac{\partial}{\partial \tau}+i \pi g V^{2}(q(\tau)+q(\theta))^{2}\right] \widehat{\rho}_{-}(\zeta, \tau, \theta)=2 \pi i \tag{6.141}
\end{equation*}
$$

The solution of this equation yields the population difference at a given time $t$, if after performing an inverse Fourier transformation one sets $\tau=\theta=t$. Of course, one can always act in this way and obtain the correct result, however we note that the same result can be obtained directly when we remember that the axillary variables $\tau$ and $\theta$ have been introduced in order to replace the time-dependent Schrödinger equation by a time-independent one with one more coordinate for the wavefunction $\psi$ or two more coordinates for the density matrix $\rho$. Therefore, (6.141) is just a representation of the equation

$$
\begin{equation*}
\left[\frac{\partial}{\partial t}+4 \pi g V^{2} q^{2}(t)\right] \rho_{-}(t)=0 \tag{6.142}
\end{equation*}
$$

with the initial condition $\widehat{\rho}_{-}(t=0)=1$ which is much easier to solve directly, without making use of such a technique. This yields

$$
\begin{equation*}
\rho_{-}(t)=e^{-4 \pi g V^{2}} \int_{0}^{t} q^{2}(x) d x, \tag{6.143}
\end{equation*}
$$

and for $q(t)=\cos [t+\varphi]$ results in

$$
\begin{equation*}
\rho_{-}(t)=\exp \left[-2 \pi g V^{2}(t+\cos (t+2 \varphi) \sin t)\right] \tag{6.144}
\end{equation*}
$$

or in dimensional variables for frequency $\omega$

$$
\begin{equation*}
\rho_{-}(t)=\exp \left[-2 \pi g V^{2}\left(t+\omega^{-1} \cos (\omega t+2 \varphi) \sin \omega t\right)\right] . \tag{6.145}
\end{equation*}
$$

This function shown in Fig. 6.14(d), demonstrates a decay with an oscillating rate.


Fig. 6.14. One band in a bi-harmonic field with frequencies $\omega_{h} \pm \omega$ comparable with the band width (a) can be considered as a two-band problem (b) with the harmonic interaction at the beat frequency $\omega$. Application of the quasienergy method (c) suggests the consideration of an equivalent multiband system with a time-independent coupling, where each harmonic of the population amplitude $\psi_{k n}$, (6.146), corresponds to an individual band. (d) The difference between the total band populations (6.145) for $W=1, \omega=5$, and the phase shifts $\phi=\pi / 4$ (dashed line) or $\phi=-\pi / 4$ (dot dash line) is a decaying function of time, with an oscillating decrement. For the field in the quantum state corresponding to a given number of photons (6.157), the dependence has no oscillations (solid line)

## Quantized Pumping Field

Thus far in our consideration, we have always assumed that the external field is classical in nature. The employed technique of introducing the additional variables $\tau$ and $\theta$ is identical to that known as the quasi-energy method briefly discussed in Sect.5.1, which is closely related to the quantum consideration of the pumping field. The main idea of the method is the following.

We substitute the anzats

$$
\begin{equation*}
\psi_{n}=\sum_{k} e^{-i k \omega t} \psi_{k n}(t) \tag{6.146}
\end{equation*}
$$

into the Schrödinger equation (6.30) with $q(t)=\cos (\omega t)$, and considering all of the coefficients corresponding to the harmonics $e^{-i k \omega t}$ arrive at the expression

$$
\begin{equation*}
i \dot{\psi}_{k n}=\left(\Delta_{n}-\omega k\right) \psi_{k n}+\frac{1}{2} \sum_{m} V_{n m}\left(\psi_{k-1, m}+\psi_{k+1, m}\right), \tag{6.147}
\end{equation*}
$$

which is known as the quasi-energy representation of (6.30) with periodic coefficients. This new equation does not contain the time-dependent coefficients and therefore can be solved by standard methods. Formally it can be considered as a multiband system as shown in Fig. 6.14(c), where each quasi-energy index $k$ corresponds to an individual band of levels enumerated either by the index $m$ or by index $n$, depending on the number of the band (1 or 2 ) to which it belongs initially.

But now we have encountered the problem of undefined initial conditions. Indeed, the requirement $\psi_{0}(t=0)=1$ with the allowance of (6.146) only states that $\sum_{k} \psi_{k 0}(0)=1$ and does not give any individual initial conditions for each amplitude $\psi_{k n}(t=0)$. This problem has deep roots in the question of the quantum statistics of the pumping field, and its solution is different for the coherent states of the quantum field corresponding to the classical limit, and for the photon number state of the field which corresponds to a definite field energy.

Let us look at (6.147) from the point of view of a quantized field. Considering the two equivalent problems of a single band in the two-frequency field at $\omega_{h} \pm \omega$ and two bands in a harmonic field at the frequency $\omega$ we favor the latter approach and denote by $\psi_{k n}$ the probability amplitude for the system to be at the level $n$ of the first band and the quantized field to be in the state with $k \gg 1$ photons at the frequency $\omega$. This corresponds to the state $|n\rangle \otimes|k\rangle$ of the compound system. By $\psi_{k m}$ we denote the analagous amplitudes for the second band, and therefore (6.147) now may be directly interpreted as the Schrödinger equation for the compound system. By introducing the operators $\widehat{D}$ and $\widehat{q}$ with the matrix elements $D_{k k^{\prime}}=\omega k \delta_{k, k^{\prime}}$ and $q_{k k^{\prime}}=\left(\delta_{k, k^{\prime}+1}+\delta_{k, k^{\prime}-1}\right) / 2$ one obtains for this system the Fourier transformed evolution operator

$$
\begin{equation*}
\widehat{U}(\varepsilon)=\left(\varepsilon-\widehat{H}_{0}+\widehat{D}-\widehat{q} \otimes \widehat{V}\right)^{-1} \tag{6.148}
\end{equation*}
$$

which actually coincides with (6.32) if we note that the operators $\widehat{D}$ and $\widehat{q}$ are just a Fourier (in $k$ ) representation of $i \frac{\partial}{\partial \tau}$ and the harmonic time-dependent profile, respectively.

There is no need to repeat here all of the derivation, since by analogy one can write equations (6.138) and (6.140) for the total populations of the bands in a form similar to (6.72)

$$
\begin{align*}
\widehat{Z}_{-}^{-1}(\zeta) \widehat{\rho}_{+}(\zeta) & =2 \pi i \\
\widehat{Z}_{+}^{-1}(\zeta) \widehat{\rho}_{-}(\zeta) & =2 \pi i \\
\widehat{Z}_{\mp}^{-1}(\zeta) & =\zeta-\widehat{D}+\widehat{D}^{\prime}+i \pi g V^{2}\left(\widehat{q} \pm \widehat{q}^{\prime}\right)^{2} \tag{6.149}
\end{align*}
$$

where the operators $\widehat{D}^{\prime}$ and $\widehat{q}^{\prime}$ act on the density operators $\widehat{\rho}_{ \pm}$from right to left. The density matrices $\widehat{\rho}_{ \pm}$do not depend on the indices enumerating the sublevels $m$ of the system, but they do depend on the photon variables $k$. Note that the first equation (6.149) after the Fourier transformation in $\zeta$ takes the form of an equation for commutators and anticommutators

$$
\begin{align*}
& i \frac{\partial}{\partial t} \widehat{\rho}_{+}=\left[\widehat{D}, \widehat{\rho}_{+}\right]-i \pi g V^{2}\left[\widehat{q},\left[\widehat{q}, \widehat{\rho}_{+}\right]\right] \\
& i \frac{\partial}{\partial t} \widehat{\rho}_{-}=\left[\widehat{D}, \widehat{\rho}_{-}\right]-i \pi g V^{2}\left\{\widehat{q},\left\{\widehat{q}, \widehat{\rho}_{-}\right\}\right\} \tag{6.150}
\end{align*}
$$

The traces of $\widehat{\rho}_{+}$and $\widehat{\rho}_{-}$will give the sum and the difference of the total band populations respectively. We concentrate on the second trace given by the second equation (6.150), since for the first one the result is evidently unity. After the trace-preserving transformation $\widehat{\rho}_{-} \rightarrow e^{-i \widehat{D} t} \widehat{\rho}_{-} e^{i \widehat{D} t}$ one obtains

$$
\begin{equation*}
\frac{\partial}{\partial t} \widehat{\rho}_{-}=-\pi g V^{2}\left\{e^{-i \widehat{D} t} \widehat{q} e^{i \widehat{D} t},\left\{e^{-i \widehat{D} t} \widehat{q} e^{i \widehat{D} t}, \widehat{\rho}_{-}\right\}\right\} \tag{6.151}
\end{equation*}
$$

which in the explicit matrix notation reads

$$
\begin{align*}
-\frac{4}{\pi g V^{2}} \frac{\partial}{\partial t} \rho_{k, k^{\prime}}= & \left(4 \rho_{k, k^{\prime}}+2 \rho_{k+1, k^{\prime}+1}+2 \rho_{k-1, k^{\prime}-1}\right) \\
& +e^{i 2 \omega t}\left(\rho_{k+2, k^{\prime}}+2 \rho_{k+1, k^{\prime}-1}+\rho_{k, k^{\prime}-2}\right) \\
& +e^{-i 2 \omega t}\left(\rho_{k-2, k^{\prime}}+2 \rho_{k-1, k^{\prime}+1}+\rho_{k, k^{\prime}+2}\right) \tag{6.152}
\end{align*}
$$

and we have omitted here the index "-". Let us now introduce $\rho_{s}=$ $\sum_{k} \rho_{k, k+s}$ and note that $\rho_{s=0}$ gives the band population difference. Equation (6.152) yields for this quantity

$$
\begin{equation*}
-\frac{1}{\pi g V^{2}} \frac{\partial}{\partial t} \rho_{s}=2 \rho_{s}+e^{i 2 \omega t} \rho_{s-2}+e^{-i 2 \omega t} \rho_{s+2} \tag{6.153}
\end{equation*}
$$

We make use of the generation function $\rho(\beta, t)=\sum_{s} \rho_{s}(t) e^{i s \beta}$ for which one immediately finds

$$
\begin{equation*}
-\frac{1}{2 \pi g V^{2}} \frac{\partial}{\partial t} \rho(\beta, t)=(1+\cos [2 \omega t+2 \beta]) \rho(\beta, t) \tag{6.154}
\end{equation*}
$$

that is

$$
\begin{equation*}
\rho(\beta, t)=C(\beta) \exp \left\{-2 \pi g V^{2}\left[t+\frac{1}{2 \omega} \sin (2 \omega t+2 \beta)\right]\right\}, \tag{6.155}
\end{equation*}
$$

and hence

$$
\begin{equation*}
\rho_{s}(t)=\int_{0}^{2 \pi} d \beta \frac{C(\beta)}{2 \pi} \exp \left\{-2 \pi g V^{2}\left[t+\frac{1}{2 \omega} \sin (2 \omega t+2 \beta)\right]-i s \beta\right\} . \tag{6.156}
\end{equation*}
$$

For the case $C(\beta)=2 \pi \delta(\beta)$ (6.156) coincides with the result (6.145) obtained for a classical field of phase $\varphi=0$. This result also follows from the consistent quantum consideration of the problem: for a quantized field in a coherent state with a large mean number of photons the initial photonic density matrix $\sim \exp \left(-i\left(k-k^{\prime}\right) \varphi\right)$ results in $\rho_{s}(t=0) \sim \exp (-i s \varphi)$, and hence according to (6.155) yields $C(\beta)=2 \pi \delta(\beta-\varphi)$, which indeed gives $C(\beta)=2 \pi \delta(\beta)$ for the field phase $\varphi=0$.

For the field initially in the state with a given number of photons $n_{0}$, the photonic density matrix has the form $\delta_{k, n_{0}} \delta_{k^{\prime}, n_{0}}$ and therefore $C(\beta)=$ $\rho(\beta, t=0)=1$. The population difference can therefore be given in terms of a modified Bessel function

$$
\begin{align*}
\rho_{-}(t)=\rho_{s=0}(t)= & \frac{e^{-2 \pi g V^{2} t}}{2 \pi} \int_{0}^{2 \pi} d \beta \exp \left[-\frac{\pi g V^{2}}{\omega} \sin (2 \omega t+2 \beta)\right] \\
& =e^{-\pi g V^{2} t} I_{0}\left(\frac{\pi g V^{2}}{\omega} \sin \omega t\right) \tag{6.157}
\end{align*}
$$

which is different from the absorption rate in a coherent field (6.145). However this expression can be obtained from (6.145) by averaging it over the field phase $\varphi$. We have thus obtained an interesting result, which shows that the complex systems absorbs radiation with the rate depending on the quantum field statistics.

### 6.5.2 Population Distribution over the Bands

We now consider the distribution of the populations over the bands. Actually, this problem has already been solved for the sum $\rho_{n+}$ of the band populations, for which the derivation coincides with that for (6.49) and (6.55) while for the population difference $\rho_{n-}$ the result is identical apart from the change of sign in the cross-term of the correlation function, such that the general expression reads

$$
\begin{equation*}
\pi g \rho_{n \pm}(t)=e^{-W \int_{0}^{t} q^{2}(x) d x} \int_{0}^{t} \cos \left(\Delta_{n} y\right) \exp \left[ \pm W \int_{0}^{t-y} q(x) q(y+x) d x\right] d y \tag{6.158}
\end{equation*}
$$

The change of sign originates from the fact that for $\rho_{n-}$ instead of the operator $\widehat{Z}$ given by (6.40), one has to make use of the operator $\widehat{Z}_{+}$given by (6.149).


Fig. 6.15. Density of the population distribution (arb. units) of the first band $g \rho_{1}=\left(\rho_{+}+\rho_{-}\right) / 2$ and of the second band $g \rho_{2}=\left(\rho_{+}-\rho_{-}\right) / 2$, as functions of time. For short times, one sees that the profiles resembling those shown in Figs .4.16, 4.21 for the level-band system and for the degenerate level-band system respectively, although in the asymptotic of long times, these profiles behave differently, showing a diffusive broadening.

As we have already seen, for small $W / \omega$ and an abrupt switch of the interaction, the population is distributed over the narrow vicinities of the resonances with Lorentzian profiles according to (6.97). This is also evidently the case for the two-band system, as long as the resonances do not overlap we have just to take into account that the even resonances correspond to the first and the odd to the second bands. We therefore concentrate here on the other extreme $W \gg \omega$ which we referred to on p. 265 and restrict ourselves to the case of a coherent field. Taking as earlier the frequency units $(\omega=1)$ and performing the integration in the exponent of (6.158) one obtains

$$
\begin{align*}
& \pi g \rho_{n \pm}(t)=e^{-W[t+\cos (t+2 \varphi) \sin t] / 2}  \tag{6.159}\\
& \int_{0}^{t} \cos \left(\Delta_{n} y\right) \exp \left[ \pm \frac{W}{2}((t-y) \cos y+\cos (t+2 \varphi) \sin (t-y))\right] d y
\end{align*}
$$

In Fig. 6.15 we show the population distribution over the first and the second bands as a function of time.

For the asymptotically long times, the population difference $\rho_{n-}$ becomes small, since the maximum index of the exponent in the integrand for the case of negative sign cannot compensate for the exponentially decreasing factor in front of the integral. The population distributions of the first and the second bands are therefore identical for long times, and in order to find them one can consider their sum $\rho_{n+}$. The main contribution comes from small $y \sim 1 / \sqrt{W}$ such that one can cast the expression in the exponent of the integrand in a Taylor series and obtain

$$
\begin{align*}
\pi g \rho_{n \pm}(t)= & \int_{0}^{\infty} \cos \left(\Delta_{n} y\right) \\
& \exp \left[\frac{-W}{2}\left(t \frac{y^{2}}{2}+y \cos (t+2 \varphi) \cos t\right)\right] d y \tag{6.160}
\end{align*}
$$

where the integration interval has been extended to infinity. Evaluation of the integral for $t \gg W$ (with the allowance of $\omega=1$ ) results in the distribution profile

$$
\begin{align*}
\pi g \rho_{+}(\Delta, t) \simeq & \sqrt{\frac{\pi}{t W}} e^{-\Delta^{2} / W t}\left\{\cos \left[\frac{\Delta}{t} \cos t \cos (t+2 \varphi)\right]\right. \\
& \left.+\operatorname{Erfi}\left(\frac{\Delta}{\sqrt{W t}}\right) \sin \left[\frac{\Delta}{t} \cos t \cos (t+2 \varphi)\right]\right\} \tag{6.161}
\end{align*}
$$

where $\operatorname{Erfi}(x)$ is the error function. The Gaussian dependence implies that $\Delta \sim \sqrt{W t}$, and therefore $\Delta / t \ll 1$, which means that in order to be consistent within the approximation we have to replace this expression by an even simpler Gaussian dependence

$$
\begin{equation*}
\pi g \rho_{+}(\Delta, t) \simeq \sqrt{\frac{\pi}{t W}} e^{-\Delta^{2} / W t} \tag{6.162}
\end{equation*}
$$

In other words the population distribution is broadening in the long-time limit, according to the diffusion law.

### 6.6 Control of Complex Quantum Systems

Thus far we have been considering the universal properties common to almost all realizations of complex quantum systems, regardless of the details of the unperturbed Hamiltonian $\widehat{H}_{0}$ and the interaction $\widehat{V}$. We now consider another aspect of such systems, namely their controllability, the possibility to make a complex system follow a predescribed evolution by choosing the strengths and types of applied external perturbations. The control goals can be of two types. One can either place the system at a certain time in a well-defined particular quantum state, begining from another well-defined initial state, or perform a well-defined transformation of an unknown initial quantum state.

In the first case we speak of quantum state manipulation, whereas the second case implies a more involved control over the unitary evolution of the entire system, given by a specific unitary transformation matrix. The latter also implies a predetermined evolution of each basis state in the Hilbert space of the system.

Though it is not evident at first glance, the justification of the statistical description of complex quantum systems and the requirement necessary for the controllability of a multilevel system have common algebraic and group theory roots. In spite of this, the ways of describing these two processes are very different. Considering the dynamics of a complex system, one ignores individual details of the spectrum and the interaction performing the ensemble average, thus assuming that the particular realization of a complex system has a small influence on the universal character of its behavior, given by a few average characteristics. For the problem of control the situation is quite the opposite: one needs to know all of the spectral and interaction details in order to achieve the desired system evolution. The techniques for solving these two problems, as well as the final results obtained, are accordingly distinct. In order to describe the universal dynamics of a complex system, one tries to find a finite algebraic expression for the ensemble averaged physical quantity of interest in terms of standard functions, and to check whenever it is possible, whether the dispersion of the results is small. The solution of the control problem is usually an algorithm which enables one to compute the control functions $f_{i}(t)$, by employing detailed information about the Hamiltonian $\widehat{H}_{0}$, the perturbations $\widehat{V}_{j}$, and the desired transformation matrix $\widehat{U}_{d}$.

The most general control problem is formulated for a Hamiltonian

$$
\begin{equation*}
\widehat{H}(t)=\widehat{H}_{0}+\sum_{j} f_{j}(t) \widehat{V}_{j} \tag{6.163}
\end{equation*}
$$

comprising a given time-independent unperturbed part $\widehat{H}_{0}$ and a given set of different perturbation operators $\widehat{V}_{j}$. Each of the perturbations may have an individual time-dependent strength $f_{j}(t)$ and the solution implies that one can specify the particular forms of $f_{j}(t)$ such that the evolution matrix

$$
\begin{align*}
& \widehat{U}(t)=\widehat{T} e^{-i} \int_{0}^{t} \widehat{H}(t) d t
\end{align*}=1-i \int_{0}^{t} \widehat{H}(t) d t+i^{2} \int_{0}^{t} \widehat{H}(t) \int_{0}^{t} \widehat{H}\left(t^{\prime}\right) d t^{\prime} d t+\ldots .
$$

assumes a predetermined form $\widehat{U}(T)=\widehat{U}_{d}$ at a given time moment $t=T$. When this is possible the system is called controllable. Here $\widehat{T}$ is the time ordering operator that has been encountered earlier in (5.2). A particular case, where all functions $f_{j}(t)$ take only two values $f_{j}(t)=0$ and $f_{j}(t)=1$, is called bang-bang control.

The nessecery condition of the system's controllability is that the Hamiltonian $\widehat{H}_{0}$ and perturbation operators $\widehat{V}_{j}$ together with their commutators
of different orders form an operator basis set which spans the entier linear space of the operators in the Hilbert space of the system. In other words, it should be possible to represent any operator in the Hilbert space of the considered quantum system as a linear combination of $\widehat{H}_{0}$, operators $\widehat{V}_{j}$, and their commutators of various orders.

We consider several aspects of the control problem, starting with the simplest examples of the two-level system. Though the consideration of these examples can be performed directly in exact analytical form, we discuss several algorithms for the numerical solution and introduce the main concepts and images convenient for the description of more sophisticated situations.

### 6.6.1 Control of Two-Level Systems

A two-level quantum system within two-dimensional Hilbert space corresponds to the Lie group $S U(2)$ and can be described with the help of the three-element algebra of Pauli matrices $\widehat{\sigma}_{x}, \widehat{\sigma}_{y}, \widehat{\sigma}_{z}$, that compose the set of generators of $S U(2)$ supplemented by the identity matrix $\widehat{I}$. This group is isomorphic to the group of three dimensional rotations, and therefore the main ideas of the control technique can be illustrated by three-dimensional images, convenient for intuitive perception.

Indeed, each Hermitian operator $\widehat{H}$ can be represented as a linear combination of Pauli matrices

$$
\begin{equation*}
\widehat{H}=X \widehat{\sigma}_{x}+Y \widehat{\sigma}_{y}+Z \widehat{\sigma}_{z}+E_{0} \widehat{I} \tag{6.165}
\end{equation*}
$$

where the constant $E_{0}$ determining the mean energy responsible for the global phase can be omitted. Therefore the vector $\bar{H}=(X, Y, Z)$ formed by three real numbers determines the Hermitian operator with zero trace. The quantum state of the two-level system can also be represented as a real threedimensional vector of unit length. Two eigenvectors of an operator $\widehat{H}$ are given by the three-dimensional vectors $(x, y, z)$ of length 1 that are parallel and antiparallel to the vector $\bar{H}$. Each state vector $(x, y, z)$ also corresponds to a unit rank density matrix

$$
\begin{equation*}
\rho=|\psi\rangle\langle\psi|=\frac{x}{2} \widehat{\sigma}_{x}+\frac{y}{2} \widehat{\sigma}_{y}+\frac{z}{2} \widehat{\sigma}_{z}+\frac{1}{2} \widehat{I} \tag{6.166}
\end{equation*}
$$

of the pure state. The coordinates $(x, y, z)$ of the state $|\psi\rangle$ and the state $|\widetilde{\psi}\rangle$ orthogonal to it, thus correspond to opposite points on the unit sphere, as shown in Fig. 6.16a). These coordinates are also known as the components of the Bloch vector.

Each unitary transformation

$$
\begin{equation*}
\widehat{U}=e^{-i t \widehat{H}}=\exp \left[-i t\left(X \widehat{\sigma}_{x}+Y \widehat{\sigma}_{y}+Z \widehat{\sigma}_{z}\right)\right] \tag{6.167}
\end{equation*}
$$

of a state vector $(x, y, z)$ under the action of a time-independent Hamiltonian is a counterclockwise rotation around the axis $\bar{H}$ by the angle $\phi=$


Fig. 6.16. Representation of the states and operators of a two-level system by three-dimensional vectors. (a) Each Hamiltonian $\hat{H}$ corresponds to a vector $\bar{H}$ with the coordinates $(X, Y, Z)$ given by the coefficients in the representation (6.165) of this Hamiltonian as the sum of Pauli matrices $\sigma_{x}, \sigma_{y}$, and $\sigma_{z}$. Two eigenstates of the Hamiltonian, with positive $(+)$ and negative $(-)$ eigenenergies, correspond to the points given by the intersection of the vector axis with the unit sphere. Coordinates $\pm(x, y, z)$ of these points are given by the components $P, Q$, and $W$ of the Bloch vector (6.166), respectively. (b) A unitary transformation $\exp (-i t \hat{H})$ can be represented by an axial vector $\bar{\phi}$, directed parallel to $\bar{H}$, which has a length equal to the rotation angle. The points representing the states move on the unit sphere. (c) Vector $\bar{U}$ representing the unitary transformation as a sum (6.168) does not conserve its length, moving under the action of a different transformation around the vector $\bar{H}$, along an elliptic trajectory.
$t\left(X^{2}+Y^{2}+Z^{2}\right)^{1 / 2}$ represented by the axial vector $\bar{\phi}=(X t, Y t, Z t)$, as shown in Fig. 6.16b). One can also derive a differential equation describing the evolution of the unitary operator $\widehat{U}$ under the action of a time-dependent Hamiltonian. Indeed, by representing the evolution operator as a linear combination

$$
\begin{equation*}
\widehat{U}=U_{0} \widehat{I}-i U_{X} \widehat{\sigma}_{x}-i U_{Y} \widehat{\sigma}_{y}-i U_{Z} \widehat{\sigma}_{z} \tag{6.168}
\end{equation*}
$$

and by calculating the products of the Pauli matrices, one can rewrite the Schrödinger equation $i \frac{\partial}{\partial t} \widehat{U}=\widehat{H} \widehat{U}$ in the form

$$
\begin{align*}
\frac{\partial}{\partial t} U_{0} & =-(\overline{H U}) \\
\frac{\partial}{\partial t} \bar{U} & =\bar{H} U_{0}+[\bar{H} \times \bar{U}] \tag{6.169}
\end{align*}
$$

employing the vector notation $\bar{U}=\left(U_{X}, U_{Y}, U_{Z}\right)$ and $\bar{H}=(X, Y, Z)$.
The unitarity of $\widehat{U}$ implies that $U_{0}^{2}+U_{X}^{2}+U_{Y}^{2}+U_{Z}^{2}=1$, which means that the length of the three-dimensional part $\bar{U}$ of the four dimensional vector $\left(U_{0}, U_{X}, U_{Y}, U_{Z}\right)$ is not conserved. Inspection of (6.169) shows that $\bar{U}$ moves along an elliptic trajectory over the surface of a cylinder directed along the vector $\bar{H}$ as shown in Fig. 6.16(c), since the component of $\bar{U}$ rotates around the axis given by the direction of $\bar{H}$ with the angular velocity $|\bar{H}|$, whereas the component $U_{\| \mid}$parallel to $\bar{H}$ oscillates with the same frequency as that of the component $U_{0}$. The latter is a manifestation of the rotation of the vector $\left(U_{\|}, U_{0}\right)$ in four-dimensional space. However, due to the unitarity condition, the three-dimensional vector $\bar{U}$ describes the evolution operator completely.

The fact that the quantum states of the system, the Hamiltonians, and the unitary transformation matrices can all be represented in the same threedimensional real vector space is a unique property of the two-level systems corresponding to the $S U(2)$ symmetry group. For other systems, the dimensionality of the real vector spaces corresponding to operators are larger as compared to the vector spaces representing the states. When we ignore the global phase of the state and allow for the normalization condition, each state of an $N$-level system given by a normalized complex-valued vector in $N$-dimensional Hilbert space corresponds to a $2 N-2$ real vector. The Hamiltonians, being Hermitian matrices, can be described by $N^{2}-1$ real parameters when we set their mean energy responsible for the global phase to zero, thus assuming that all Hamiltonians have zero trace. This is also valid for the unitary matrices generated by such Hamiltonians via the relation $\widehat{U}=\exp (-i t \widehat{H})$.

## Iterative Transformations of States

The control of two-level systems is not a challenging task, and usually solutions can be found with the help of relatively simple algebraic calculations. However, it is useful to consider several control algorithms as operations in the vector space, since this illustrates the main principle of finding the control amplitudes $f_{j}(t)$ for large systems. We mainly restrict ourselves to the so-called "bang-bang" control when the amplitudes assume only the values 0 or 1 during certain time intervals.

Let us assume that $\widehat{H}_{0}=0$ and that we have only two different perturbations $\widehat{V}_{1}=\widehat{\sigma}_{x}$ and $\widehat{V}_{2}=\widehat{\sigma}_{z}$ at our disposal. The commutator $\left[\widehat{V}_{1}, \widehat{V}_{2}\right]=2 i \widehat{\sigma}_{y}$ provides us with the missing third operator and the operator basis becomes complete and the system is therefore controllable. For this simplest of examples, we will illustrate the geometrical meaning of the algorithms that allow one to control the system.

The simplest iterative algorithm allows one to approach a given quantum state $|f\rangle$ starting from an initial state $|i\rangle$. We apply to the system initially in the state $|i\rangle$, the perturbations $\widehat{V}_{1}$ and $\widehat{V}_{2}$ in sequence, trying to approach
the desired state $|f\rangle$ with each iteration step. As a result, the initial state experiences the series of transformations $\left|i_{N}\right\rangle=\exp \left(-i t_{N} \widehat{V}_{1,2}\right)\left|i_{N-1}\right\rangle$ where $\widehat{V}_{1,2}$ can be either $\widehat{V}_{1}$ or $\widehat{V}_{2}$, depending on the parity of $N$.

At each step of the iteration we try to maximize the scalar product

$$
\begin{equation*}
\left\langle f \mid i_{N}\right\rangle=\langle f| \exp \left(-i t_{N} \widehat{V}_{1,2}\right)\left|i_{N-1}\right\rangle \tag{6.170}
\end{equation*}
$$

by the proper choice of the parameter $t_{N}$ denoting the exposure time to the perturbation. After $N$ iterations we therefore achieve the state

$$
\begin{equation*}
\left|i_{N}\right\rangle=\widehat{U}|i\rangle \tag{6.171}
\end{equation*}
$$

which approaches the desired state $|f\rangle$ as a result of the application of the iteratively constructed operator

$$
\begin{equation*}
\widehat{U}=\exp \left(-i t_{N} \widehat{V}_{1,2}\right) \ldots \exp \left(-i t_{3} \widehat{V}_{1}\right) \exp \left(-i t_{2} \widehat{V}_{2}\right) \exp \left(-i t_{1} \widehat{V}_{1}\right)|i\rangle . \tag{6.172}
\end{equation*}
$$

as illustrated in Fig. 6.17(a).


Fig. 6.17. The iterative procedure compared to the direct transition. (a) By successive application of two different transformations, $\exp \left(-i \hat{V}_{1} t\right)$ and $\exp \left(-i \hat{V}_{2} t\right)$, suggested by the two Hamiltonians $\hat{V}_{1}=\hat{\sigma}_{x}$ and $\hat{V}_{2}=\hat{\sigma}_{z}$ one can reach the final state $|f\rangle$ from the initial state $|i\rangle$ by choosing the variable $t_{N}$ at each step in such a way that the scalar product of the transformed state and the target state is maximum. (b) The same transformation can be performed in two steps only by finding the maximum scalar product as a function of two variables, $t_{1}$ and $t_{2}$. However this approach may not result in the coincidence of the transformed state with the target state and can result only in the closest approach if the order of the applied transformations is wrong, as shown by the spots with the dashed line contour. The scalar product then has a local maximum, different from the global maximum value 1.

The iterative procedure is not a very efficient way to perform the desired transformation. Indeed, it depends on a large number of parameters $t_{N}$, whereas only two parameters are required for determining the position of a point on the unit sphere representing a state of the two-level system.

Therefore the objective can be achieved by application of only two transformations

$$
\begin{equation*}
\left|i_{2}\right\rangle=\exp \left(-i T_{2} \widehat{V}_{2}\right) \exp \left(-i T_{1} \widehat{V}_{1}\right)|i\rangle \tag{6.173}
\end{equation*}
$$

However, for this purpose, one needs to perform a simultaneous search for the parameters $T_{1}$ and $T_{2}$ that yield the maximum scalar product $\left\langle f \mid i_{2}\right\rangle$, as illustrated in Fig. 6.17(b). The main obstacle in this approach is that the problem may not always have a solution for a given transformation sequence. The natural way to overcome this obstacle is to add an extra transformation

$$
\begin{equation*}
\left|i_{3}\right\rangle=\exp \left(-i T_{3} \widehat{V}_{1}\right) \exp \left(-i T_{2} \widehat{V}_{2}\right) \exp \left(-i T_{1} \widehat{V}_{1}\right)|i\rangle \tag{6.174}
\end{equation*}
$$

and look for the maximum in the space of three variables instead of two. The solution of the latter problem is not unique, and the set of solutions form a one-dimensional curve in the three-dimensional space of the parameters $\left(T_{1}, T_{2}, T_{3}\right)$.

## Iterative Construction of a Unitary Transformation

Prior to the consideration of an algorithm allowing one to construct a desired unitary transformation, we take advantage of the simplicity of the two-level system and solve the problem directly by calculation of the most direct transformation. For the case $\widehat{V}_{1}=\widehat{\sigma}_{x}$ and $\widehat{V}_{2}=\widehat{\sigma}_{z}$ under consideration, the direct multiplication of matrices yields

$$
\begin{align*}
\widehat{U}= & \exp \left(-i T_{3} \widehat{V}_{1}\right) \exp \left(-i T_{2} \widehat{V}_{2}\right) \exp \left(-i T_{1} \widehat{V}_{1}\right) \\
= & \cos \left(T_{3}+T_{1}\right) \cos \left(T_{2}\right) \widehat{I}-i \sin \left(T_{3}+T_{1}\right) \cos \left(T_{2}\right) \widehat{\sigma}_{x} \\
& +i \sin \left(T_{3}-T_{1}\right) \sin \left(T_{2}\right) \widehat{\sigma}_{y}-i \cos \left(T_{3}-T_{1}\right) \sin \left(T_{2}\right) \widehat{\sigma}_{z} . \tag{6.175}
\end{align*}
$$

Comparison of (6.175) and (6.168) allows us to find the vector $\left(U_{9} U_{X}, U_{Y}, U_{Z}\right)$. The solution of the equation $\widehat{U}\left(T_{1}, T_{2}, T_{3}\right)=\widehat{U}_{d}$ can be found for any desired unitary matrix $\widehat{U}_{d}=U_{0} \widehat{I}-i U_{X} \widehat{\sigma}_{x}-i U_{Y} \widehat{\sigma}_{y}-i U_{Z} \widehat{\sigma}_{z}$, and reads

$$
\begin{align*}
& T_{1}=\frac{1}{2} \arcsin \frac{U_{Y}}{\sqrt{U_{X}^{2}+U_{Y}^{2}}}+\frac{1}{2} \arcsin \frac{U_{Z}}{\sqrt{U_{0}^{2}+U_{Z}^{2}}} \\
& T_{2}=\arccos \left(-\sqrt{U_{X}^{2}+U_{Y}^{2}}\right) \\
& T_{3}=\frac{1}{2} \arccos \frac{-U_{Y}}{\sqrt{U_{X}^{2}+U_{Y}^{2}}}+\frac{1}{2} \arccos \frac{U_{Z}}{\sqrt{U_{0}^{2}+U_{Z}^{2}}} . \tag{6.176}
\end{align*}
$$

One can verify directly that the coefficients $U$ given by the trigonometric expressions do indeed satisfy the requirement $U_{0}^{2}+U_{X}^{2}+U_{Y}^{2}+U_{Z}^{2}=1$.

The solution (6.176) can also be found numerically by minimizing the function

$$
\begin{equation*}
\operatorname{Tr}\left[\left(\widehat{U}\left(T_{1}, T_{2}, T_{3}\right)-\widehat{U}_{d}\right)\left(\widehat{U}\left(T_{1}, T_{2}, T_{3}\right)-\widehat{U}_{d}\right)^{*}\right] \tag{6.177}
\end{equation*}
$$

which gives the square of the distance between two matrices vanishing for the desired solution. However, if by analogy to the state control we try to find a solution minimizing successively the function

$$
\begin{equation*}
\operatorname{Tr}\left\{\left[\exp \left(-i t_{N} \widehat{V}_{1,2}\right) \ldots \exp \left(-i t_{2} \widehat{V}_{2}\right) \exp \left(-i t_{1} \widehat{V}_{1}\right)-\widehat{U}_{d}\right] \times c . c\right\} \tag{6.178}
\end{equation*}
$$

for each new added operator $\exp \left(-i t_{N} \widehat{V}_{1,2}\right)$, the procedure may converge to a local minimum different from zero and not give the desired result. We


Fig. 6.18. The iterative procedure for a predescribed unitary transformation. The trajectories of the vector (6.168), changing with the course of the transformations, are shown by the long dashed and dotted arrows for the direct and iterative transformations, respectively. (a) The desired transformation represented by the vector $\bar{U}$ (bold solid arrow) can be achieved directly by three successive transformations that are found by minimizing the functional (6.176). The iterative approach (6.177) may result in a transformation $\bar{U}_{i}$ which cannot be improved, since the next iteration corresponds to displacement in the direction perpendicular to that required. (b) The situation can be improved by linear iterations based on the linear expansion $\bar{T}=\bar{T}^{(0)}+\bar{\tau}$ of the operator (6.171) near the point $\left(T_{1}^{(0)}, T_{2}^{(0)}, T_{3}^{(0)}, T_{4}^{(0)}\right)$ corresponding to $\bar{U}_{i}$. Linear corrections (dash-dot arrows) allow one to iteratively achieve the solution $\bar{U}$.
illustrate this in Fig. 6.18(a) showing the situation where the last added operator moves the vector $\bar{U}$ in a direction perpendicular to that going towards the solution. Note that a similar complication resulting from the presence of local minima for the function (6.177) can also be encountered for systems with a large Hilbert space dimension.

If the iteration procedure given by the function (6.178) or any other iterative algorithm based on a incomplete set of the variables $\{t\}$ has stopped close
enough to the global minimum of the function (6.177), the exact solution can still be found by successive application of the linear approximation which involves all variables. Though for low-dimensional systems this technique does not give a substantial advantage as compared to the direct minimization of (6.177), for large systems it is much more convenient and requires much less numerical work.

Let us consider this approach in more detail. For a set of variables $\left\{T_{n}^{(0)}\right\}$ close to the solution $\left\{T_{n}\right\}$, the difference

$$
\begin{equation*}
\exp \left(-i T_{N}^{(0)} \widehat{V}_{1,2}\right) \ldots \exp \left(-i T_{2}^{(0)} \widehat{V}_{2}\right) \exp \left(-i T_{1}^{(0)} \widehat{V}_{1}\right)-\widehat{U}_{d} \tag{6.179}
\end{equation*}
$$

can be expanded into a Taylor series over the variations $\left\{\delta T_{n}\right\}$ of the variables $\left\{T_{n}\right\}$ up to the first order near the point $\left\{T_{n}^{(0)}\right\}$, that is

$$
\begin{align*}
\widehat{U}_{\delta T}= & \exp \left(-i T_{N} \widehat{V}_{1,2}\right) \ldots \exp \left(-i T_{2} \widehat{V}_{2}\right) \exp \left(-i T_{1} \widehat{V}_{1}\right)-\widehat{U}_{d}  \tag{6.180}\\
& +\sum_{n}^{N} \delta T_{n} \frac{\partial}{\partial T_{n}^{(0)}} \exp \left(-i T_{N}^{(0)} \widehat{V}_{1,2}\right) \ldots \exp \left(-i T_{2}^{(0)} \widehat{V}_{2}\right) \exp \left(-i T_{1}^{(0)} \widehat{V}_{1}\right)
\end{align*}
$$

Calculation of the derivatives is the most time-consuming part of the numerical work but must be done only once. With the help of these derivatives, we construct an algorithm which allows one to find the exact solution $T_{n}=T_{n}^{(0)}+\delta T_{n}$ as a result of the iterations over $\delta T_{n}$. Once the derivatives are found, one can find the first approximation of the variables $\delta T_{n}$ solving the system of linear equations $\widehat{U}_{\delta T}=0$ with $\widehat{U}_{\delta T}$ given by (6.180).

For the two-level system under consideration, the matrix equations $\widehat{U}_{\delta T}=$ 0 formally correspond to a system of four linear equations. However due to the unitarity of $\widehat{U}$, the actual number of linearly independent equations is three, as can be shown when we employ the Pauli matrices as the basis set in the matrix space according to (6.168) and take into account the unitarity requirement $U_{0}^{2}+U_{X}^{2}+U_{Y}^{2}+U_{Z}^{2}=1$. Hence for $N>3$ the number of equations is less than the number of variables. In this representation for $N=4$ the equation $\widehat{U}_{\delta T}=0$ adopts the matrix form

$$
\left(\begin{array}{llll}
M_{X, 1} & M_{X, 2} & M_{X, 3} & M_{X, 4}  \tag{6.181}\\
M_{Y, 1} & M_{Y, 2} & M_{Y, 3} & M_{Y, 1} \\
M_{Z, 1} & M_{Z, 2} & M_{Z, 3} & M_{Z, 3}
\end{array}\right)\left(\begin{array}{c}
\delta T_{1} \\
\delta T_{2} \\
\delta T_{3} \\
\delta T_{4}
\end{array}\right)=\left(\begin{array}{c}
\delta U_{X} \\
\delta U_{Y} \\
\delta U_{Z}
\end{array}\right)
$$

with

$$
\begin{align*}
M_{j n} & =\frac{\partial}{\partial T_{n}^{(0)}} \operatorname{Tr}\left[\widehat{\sigma}_{j} e^{-i T_{4}^{(0)} \widehat{V}_{2}} e^{-i T_{3}^{(0)} \widehat{V}_{1}} e^{-i T_{2}^{(0)} \widehat{V}_{2}} e^{-i T_{1}^{(0)} \widehat{V}_{1}}\right] \\
\delta U_{j} & =\operatorname{Tr}\left[\widehat{\sigma}_{j} e^{-i T_{4} \widehat{V}_{2}} e^{-i T_{3} \widehat{V}_{1}} e^{-i T_{2} \widehat{V}_{2}} e^{-i T_{1} \widehat{V}_{1}}-\widehat{\sigma}_{j} \widehat{U}_{d}\right] \tag{6.182}
\end{align*}
$$

where the Pauli matrices have been employed for a basis of the operator space.

As the first step we solve (6.181) with respect to $\delta T_{1}, \delta T_{2}$, and $\delta T_{3}$ since the fourth column of the matrix is perpendicular to the vector $\delta U$, which was precisely the reason for the inefficiency of the iterative variations (6.178). Then we substitute the new values of $T_{n}=T_{n}^{(0)}+\delta T_{n}$ into the second equation (6.182), find new $\delta \bar{U}$ and repeat the iteration. We note that the deviations $\delta \bar{U}$ have changed as a result of the iteration, whereas the matrix $\widehat{M}$ is retained in its original form. This is the main difference between the employed method and the steepest descent technique which requires recalculation of all of the derivatives at each step of the iteration. Nevertheless, if we are close enough to the exact solution, at each step of the iteration the displacement of variables $\delta T_{n}$ occurs at a small angle to the direction of the steepest descent and hence this procedure will converge, since for $\delta \bar{U} \rightarrow 0$ the parameters $\delta T_{n} \rightarrow 0$ as well. We illustrate this in Fig. 6.18(b).

## Control by a Periodic Field

One can control the two-level system completely by applying a harmonic external field of a given polarization, frequency and amplitude. These three parameters determine the three parameters of the effective Hamiltonian. For an external field $V_{x}=\mathcal{E}_{x} d / 2$ polarized in the $x$-direction we have obtained (3.9) describing the evolution of the state vector $\left(\psi_{0}, \psi_{1}\right)$ in the resonant approximation. For an arbitrary polarization, this equation adopts the form

$$
i \frac{\partial}{\partial t}\binom{\psi_{1}}{\psi_{0}}=\left(\begin{array}{cc}
\frac{\Delta}{2} & V_{x}-i V_{y}  \tag{6.183}\\
V_{x}+i V_{y} & -\frac{\Delta}{2}
\end{array}\right)\binom{\psi_{1}}{\psi_{0}}
$$

where $\Delta=E_{1}-E_{0}-\omega$ is the detuning of the field quantum energy from the energy of the transition. Here we do not include the initial conditions in the form of $\delta$-functions and shift the energy scale such that the mean value of the interaction vanishes.

The Hamiltonian of the system has an evident representation (6.165)

$$
\widehat{H}=\frac{\Delta}{2} \widehat{\sigma}_{z}+V_{x} \widehat{\sigma}_{x}+V_{y} \widehat{\sigma}_{y}
$$

with $X=V_{x}, Y=V_{y}$, and $Z=\Delta / 2$. Therefore any unitary evolution $\widehat{U}=\exp (-i t \widehat{H})$ can be constructed directly by choosing the strength and the orientation of the external field polarization along with the detuning. The only restriction of this method is the applicability of the resonant approximation, which requires that the Rabi frequency $\Omega=\left(\Delta^{2}+\mathcal{E}^{2} d^{2}\right)^{1 / 2}$ should be much smaller than the field frequency $\omega$. This implies that the detuning and the interactions should be small, whereas the desired transformation can be achieved on a time-scale much longer than the field period.

This type of control can also be employed for larger systems, provided they satisfy some requirements that are in close relation with the general algebraic requirements necessary for the controllability of generic quantum systems.

### 6.6.2 Holonom and Non-Holonom Systems

The universal character of complex quantum system dynamics relies on the fact that it is possible for the generic system with an unperturbed Hamiltonian $\widehat{H}_{0}$ and a perturbation $V$ to occupy any of the levels close in energy to the initially populated state, while the state vector of the system, evolving with the course of time, can approach any linear combination of the energy eigenfunctions corresponding to these levels. Mathematically this implies the validity of the bracket generation condition: for almost any pair $\widehat{V}$ and $\widehat{H}_{0}$ the set of all possible commutators $\left[\widehat{V},\left[\widehat{V},\left[\ldots,\left[\widehat{H}_{0}\right]\right] \ldots\right]\right]$, $\left[\widehat{H}_{0},\left[\widehat{H}_{0},[\ldots,[\widehat{V}]] \ldots\right]\right], \ldots$ spans the entire $N^{2}$-dimensional space of the operators, which is restricted only by the condition of Hermiticity and antiHermiticity for the even and odd orders of the commutation, respectively. Roughly speaking, the bracket generation condition is necessary for the complexity of a system and it becomes sufficient condition when all the linearly independent commutators are of the same order of magnitude. If the latter is not the case, the universality manifests itself only on the energy scale corresponding to the smallest typical size of the commutators and in the time domain given by the inverse of this value. At larger energy scales, or for shorter time intervals, the system may manifest individual behavior.

In special cases, when all of the commutators are restricted to a lowerdimensional subspace and form there a closed subalgebra, with a number of elements smaller than $N^{2}$, the system may exhibit individual behavior at any energy and at all time-scales. In the extreme case of a low-dimensional subalgebra, the multilevel system may have an exact analytical solution, thus being equivalent to a simpler quantum system. The two-level system with the three-element algebra $s u(2)$ of the Pauli matrices $\sigma_{x}, \sigma_{y}, \sigma_{z}$, along with the harmonic oscillator with the closed three-dimensional su(1,1) algebra of $a \dagger$, $a, a a \dagger$ are two examples of such simple systems.

The bracket generation condition is also necessary for complete control over the dynamics of quantum systems. Indeed, considering the timedependent systems in Sect. 5.1, we have seen that even in the simplest case of the evolution $\widehat{U}=e^{-i t_{B} \widehat{B}} e^{-i t_{A} \widehat{A}}$ under the action of two Hamiltonians $\widehat{A}$ and $\widehat{B}$ applied in sequence during the time intervals $t_{A}$ and $t_{B}$, respectively, all series of different order commutators appeared in the Baker-CampbellHausdorff formula (5.4). The coefficients with which these commutators enter the linear combination in the exponent, were powers of the exposition times $t_{A}$ and $t_{B}$. For a time dependent perturbation $f(t) \widehat{V}$ acting together with the unperturbed Hamiltonian $\widehat{H}_{0}$ the operator structure of the evolution matrix
is also of the type suggested by the Baker-Campbell-Hausdorff formula: it is an exponent of a linear combination of different commutators

$$
\begin{align*}
\widehat{U}(t)= & \exp \left\{-i \widehat{H}_{0} t-i \sum_{j} \widehat{V}_{j} \int_{0}^{t} f_{j}(t) d t+\sum_{j} c_{j 0}\left[f_{i}(t)\right]\left[\widehat{V}_{j}, \widehat{H}_{0}\right]\right. \\
& +\sum_{j k} c_{k j 0}\left[f_{i}(t)\right]\left[\widehat{V}_{k},\left[\widehat{V}_{j}, \widehat{H}_{0}\right]\right]+\ldots \\
& \left.+\sum_{\ldots j} c_{00 \ldots j}\left[f_{i}(t)\right]\left[\widehat{H}_{0},\left[\widehat{H}_{0},\left[\ldots, \widehat{V}_{j}\right] \ldots\right]\right]+\ldots\right\} \tag{6.184}
\end{align*}
$$

However, in this case, the corresponding coefficients $c$ are complex functionals of all functions $f_{j}(t)$. If the commutators of $\widehat{V}_{j}$ and $\widehat{H}_{0}$ span all $N^{2}$ dimensional operator space corresponding to the $N$-dimensional Hilbert space of the system, one can try to find the functions $f_{j}(t)$ that give the desired evolution of the system. Usually this requires much more effort since the functional relations between set of the coefficients $c$... and the set of the functions $f_{j}(t)$ is complicated and in most cases is unknown explicitly.

The bracket generation condition is equivalent to the property known in analytical mechanics as holonomicity. The essence of this concept is the following. The infinitesimal evolution of a mechanical system given by a set of local differential Hamilton equations relates the coordinates and momenta of the system in two sequential moments of time. One can often find integrals of these equations that are quantities conserved throughout the evolution. For a conservative system, the most well-known integral is the total energy, whereas the energy conservation law is the global constraint following from the local equations of motion. A local constraint, resulting in a global one, is designated as integrable or holonom. If the number of holonom constraints equals the number of degrees of freedom, the system is completely integrable or completely holonom. The motion of the completly holonom system can be separated to independent motions of differend degrees of freedom. In this case, a set of coordinates can be introduced in such a way that the motion of the system is equivalent to an independent evolution of different coordinates, and the motion of a completely holonom mechanical system with $n$ degrees of freedom occurs along the surface of an $n$-dimensional torus in $2 n$-dimensional phase space. A system is called partially holonom when the number of holonom constraints is less than the number of degrees of freedom, and for no constraints the system is non-holonom. The evolution of a non-holonom system occurs in entire (unconstrained) phase space.

There is a simple, sufficient criterion ensuring that the bracket generation condition holds for a pair of operators $\widehat{A}$ and $\widehat{B}$. Each of the operators should have non-degenerate eigenvalues $A_{j} \neq A_{k}$ with all frequencies different $A_{j}-A_{l} \neq A_{k}-A_{s}$, and in the representation where one operator, say $\widehat{A}$, is diagonal, the other one should not have zero matrix elements $B_{j, k} \neq 0$.

Note that this very condition allows one to construct any desired evolution with the Hamiltonian $\widehat{H}$ by applying harmonic external fields in complete analogy with the case of a two-level system (6.183). In fact, considering each pair of eigenvectors of $\widehat{A}$ as an independent two-level system, and applying periodic fields at the frequencies $\omega_{j, k}$ detuned from the corresponding transition frequency $A_{j}-A_{k}$ by the values $H_{j, j}-H_{k, k}$ and with the amplitudes $v_{j, k}=H_{j, k} / B_{j, k}$ one can construct any desired $\widehat{H}$ in the rotating wave approximation.

## Partially Controlled Systems

We now demonstrate an analogy between the holonom dynamics of a mechanical system and the uncontrollable dynamics of the simplest compound quantum system comprising a pair of two-level systems. At this example we draw a parallel between non-controllability and holonomicity by introducing quantum variables that behave like classical coordinates of holonom mechanical systems moving along tori surfaces. We start with the simplest case of two non-interacting two-level parts, each of which can be exposed to the action of just one Hamiltonian, $\widehat{\sigma}_{x}$ for the first part and $\widehat{\Sigma}_{z}$ for the second part. We will employ capital letters to designate the Pauli operators of the second two-level system. The system has a four-dimensional Hilbert space spanned by the orthogonal states $|0,0\rangle,|0,1\rangle,|1,0\rangle$, and $|1,1\rangle$, whereas the space of Hermitian operators consists of 16 elements, including the identity operator. The operators $\widehat{I}, \widehat{\sigma}_{x}, \widehat{\sigma}_{y}, \widehat{\sigma}_{z}, \widehat{\Sigma}_{x}, \widehat{\Sigma}_{y}, \widehat{\Sigma}_{z}$ along with nine binary direct products $\widehat{\sigma}_{j} \otimes \widehat{\Sigma}_{k}$ with $j, k=x, y, z$ can be selected as the basis spanning the operator space.

This system is apparently non-controllable since the algebra consists of only two operators $\widehat{\sigma}_{x}$ and $\widehat{\Sigma}_{z}$ that moreover commute corresponding to different Hilbert subspaces. Therefore, one cannot create new linearly independent operators via bracket generation, and the creation of controllability cannot be satisfied, since two operators $\widehat{\sigma}_{x}$ and $\widehat{\Sigma}_{z}$ alone cannot span the entire 16 -dimensional operator space. Furthermore, neither of the sub-systems is controllable individually, since the motion of the Bloch vectors representing these subsystems is just a simple rotation around the $x$-axis for the first part and around the $z$-axis for the second part of the system, that do not span the corresponding Bloch spheres. Let us construct an adequate image corresponding to this motion by representing the quantum state geometrically.

A complex vector in four-dimensional Hilbert space depends on six real parameters, provided we ignore the global phase of the state vector and allow for the normalization condition. Therefore it is impossible to depict this object as a single vector in three-dimensional space, and at least a pair of such vectors is required. When two parts of the systems are completely independent, the state vector is a direct product of the state vectors of the parts $|\Psi\rangle=\left|\psi_{1}\right\rangle\left|\psi_{2}\right\rangle$. Each part then has its own Bloch vector of unit length. By
simply appending one of the vectors to the other, scaled by a certain factor for convenience of presentation, we can construct a composition representing a quantum state of the system. The sum of the Bloch vectors contains only part of the information about the state, in contrast to the composition. Taken just by itself, the sum is insufficient for a complete description, and supplementary parameters of the composition are required for this purpose. We specify these parameters by considering different trajectories of the sum.

The initial condition $\left|\psi_{1}\right\rangle=|1\rangle,\left|\psi_{2}\right\rangle=(|1\rangle+|0\rangle) / \sqrt{2}$ corresponds to the eigenstates of the operators $\widehat{\sigma}_{z}$ and $\widehat{\Sigma}_{x}$, respectively, that are given by the vectors $(0,0,1)$ and $(1,0,0)$ on the Bloch spheres related to the first and second parts. Under the action of the Hamiltonians $\widehat{\sigma}_{x}$ and $\widehat{\Sigma}_{z}$ these vectors start to rotate in the $(z, y)$ and in the $(x, y)$ planes, respectively, and the entire evolution can be represented as a motion on the surface of a torus, as shown in Fig. 6.19 Note that each torus corresponds to a certain initial condition, which is not necessarily of the type $\left|\psi_{1}\right\rangle\left|\psi_{2}\right\rangle$, but can also include the entangled states $\alpha\left|\psi_{1}\right\rangle\left|\psi_{2}\right\rangle+\beta\left|\widetilde{\psi}_{1}\right\rangle\left|\widetilde{\psi}_{2}\right\rangle$ with $\left\langle\psi_{i} \mid \widetilde{\psi}_{i}\right\rangle=0$. Whatever the strengths $f(t)$ and $g(t)$ of the Hamiltonian $g(t) \widehat{\sigma}_{x}+f(t) \widehat{\Sigma}_{z}$ are, the evolution cannot bring the system off the surface of the torus, which is the holonom constraint of the system motion. Therefore in order to specify the state of the system, one has to indicate not only the position of the composition of the Bloch vectors, but also to specify the parameters of the torus to which it belongs. The parameters of tori, namely the vertical position of the equatorial plane, the larger radius, the eccentricity, and the smallest radius are independent integrals of motion that specify the holonom constraint as shown in Fig. 6.19(b), (c), (d).

Let us consider a more involved example of an uncontrollable system. We take two Hamiltonians $\widehat{\sigma}_{x}$ and $\widehat{\Sigma}_{z} \otimes \widehat{\sigma}_{z}$ and find that together with their commutator $\widehat{\Sigma}_{z} \otimes \widehat{\sigma}_{y}$ they form a closed, three-element algebra and together with the identity matrix span a four-dimensional subspace of the entire 16dimensional space of all $4 \times 4$ Hermitian matrices. The latter also includes the operators $\widehat{\sigma}_{y}, \widehat{\sigma}_{z}, \widehat{\Sigma}_{x}, \widehat{\Sigma}_{y}, \widehat{\Sigma}_{z}, \widehat{\Sigma}_{z} \otimes \widehat{\sigma}_{x}, \widehat{\Sigma}_{y} \otimes \widehat{\sigma}_{j}, \widehat{\Sigma}_{z} \otimes \widehat{\sigma}_{j},(j=x, y, z)$ that can never be obtained by the commutation of $\widehat{\sigma}_{z}$ and $\widehat{\Sigma}_{z} \otimes \widehat{\sigma}_{x}$.

Under the action of the Hamiltonian $\widehat{\Sigma}_{z} \otimes \widehat{\sigma}_{x}$, the trajectory of the composition of Bloch vectors moves along a torus surface in such a way that a $2 \pi$ rotation of the large radius corresponds to same rotation of the small one. After the application of the Hamiltonian $\widehat{\sigma}_{z}$ the composition moves to another torus with the same larger radius and the same position of the equatorial plane, but with a different small radius, as illustrated in Fig. 6.20(a). Having returned to the original torus by another application of the Hamiltonian $\widehat{\sigma}_{z}$, the trajectory can never return precisely to its initial position, but acquires a shift. Therefore all of the inner part of the largest torus for a given larger radius and equatorial plane can be filled by the trajectories when we apply two Hamiltonians in sequence. In other words, the set of states accessible from


Fig. 6.19. A pair of two-level systems moving independently under the action of a single operator for each system is not completely controllable. The evolution of the pair, considered as a single uncontrollable compound quantum system, can be represented as a motion along the surface of a torus with the parameters depending on the initial condition. The torus is formed by the rotation of the Bloch vector $\bar{b}$ of the first part around the end of the Bloch vector $\bar{B}$ of the second part, whereas the plane $(\bar{Z}, \bar{B})$ is considered as the $(z y)$ plane of the first part. This image is valid even for the case when the Bloch vectors of the second (b) or the first (c) are not in the planes $(z, y)$ and $(x, y)$ respectively. The image can also be extended to the case of entangled states that cannot be factorized as a direct product. In this case, $\bar{B}$ is not in the unit sphere, and has a length equal to the population difference of the factorizable components. The phase difference can also be visualized by the ellipticity of the tori cross-sections or by inclination of the equatorial plane with respect to the $Z$ axis.
an initial state by control over the operators $\widehat{\sigma}_{x}$ and $\widehat{\Sigma}_{z} \otimes \widehat{\sigma}_{z}$ comprising a three-element subalgebra is a three-dimensional variety. By employing more complicated Hamiltonians, one can obtain subalgebras with more elements which imposes fewer holonom constraints on the system dynamics and results in a variety of accessible states of a larger spatial dimensionality.

Now let us take the following Hamiltonians

$$
\begin{align*}
& \widehat{A}=\widehat{\sigma}_{z}+\widehat{\Sigma}_{x}+\widehat{\Sigma}_{x} \otimes \widehat{\sigma}_{x}+\widehat{\Sigma}_{z} \otimes \widehat{\sigma}_{y} \\
& \widehat{B}=\widehat{\sigma}_{x}+\widehat{\sigma}_{z}+\widehat{\Sigma}_{y}+\widehat{\Sigma}_{x} \otimes \widehat{\sigma}_{x}+\widehat{\Sigma}_{y} \otimes \widehat{\sigma}_{x} \tag{6.185}
\end{align*}
$$



Fig. 6.20. (a) Evolution of a couple of two-level systems under the action of two Hamiltonians that belong to a three-element subalgebra. The composition $\bar{b}$ of the Bloch vectors may now go from one torus to another with the same axis (the big radius, and the equatorial plane are identical). Holonom constraints have now only three parameters, while the constrained motion occupies three-dimensional sets, that is inner parts of the tori of smaller radius 1 and different positions of the equatorial planes. (b) The spectrum of the Hamiltonian $\hat{H}_{0}$ and the evolution operator over a period $\hat{U}=\exp (-i \hat{H} T)$. Static perturbation $V$ repels the levels close in energy. A harmonic of a time-dependent perturbation $V_{\omega}$ repels the levels separated by the energy of the corresponding quantum and leads to a more uniform distribution of the eigenstates of $\hat{U}$ along the unit circle.
chosen more or less randomly among the other combinations of a rather general structure. Their commutator gives a new, linearly independent combination which does not commute with the Hamiltonians yielding two more linearly independent combinations, and so on, until after eight cycles of commutation the number of all possible linearly independent combinations reaches the maximum possible number 15. This is an example of non-holonom system, which does not have constraints and can be controlled completely. The trajectory of the sum of the Bloch vectors can occupy the entire space restricted by the normalization requirement limiting the maximum length, and moreover each of two Bloch vectors comprising the composition occupies a three-dimensional set, and thus the state vector moves in a sixt-dimensional space.

## Control of Non-Holonom Systems

Let us try to construct a bang-bang control for the two interacting twolevel systems with the Hamiltonians (6.185) assuming that the interaction $\widehat{H}_{0}=\widehat{\sigma}_{z}+\widehat{\Sigma}_{x} \otimes \widehat{\sigma}_{x}$ is always present, whereas two different perturbations $\widehat{V}_{A}=\widehat{\Sigma}_{x}+\widehat{\Sigma}_{z} \otimes \widehat{\sigma}_{y}$ and $\widehat{V}_{B}=\widehat{\sigma}_{x}+\widehat{\Sigma}_{y}+\widehat{\Sigma}_{y} \otimes \widehat{\sigma}_{x}$ are switched on and off in sequence with the time-independent amplitudes $f_{A, B}(t)=1$ during the
exposition time intervals $T_{1}, T_{2}, \ldots, T_{16}$. Here, the even subscripts correspond to the perturbation $\widehat{V}_{B}$ with $f_{A}(t)=0, f_{B}(t)=1$ and the odd to $\widehat{V}_{A}$ with $f_{B}(t)=0, f_{A}(t)=1$. Evidently $\int f_{A}(t) f_{B}(t) d t=0$ and the total number for the switching required for complete control in the case of four-dimensional Hilbert space amounts either to 16 , or to 15 if we ignore the global phase. By analogy to a single two-level system with the evolution operator (6.175) we write

$$
\begin{align*}
& \widehat{U}\left(\left\{T_{i}\right\}\right)=e^{-i T_{16} \widehat{B}} e^{-i T_{15} \widehat{A}} e^{-i T_{14} \widehat{B}} e^{-i T_{13} \widehat{A}} e^{-i T_{12} \widehat{B}} e^{-i T_{11} \widehat{A}} e^{-i T_{10} \widehat{B}} \\
& e^{-i T_{9} \widehat{A}} e^{-i T_{8} \widehat{B}} e^{-i T_{7} \widehat{A}} e^{-i T_{6} \widehat{B}} e^{-i T_{5} \widehat{A}} e^{-i T_{4} \widehat{B}} e^{-i T_{3} \widehat{A}} e^{-i T_{2} \widehat{B}} e^{-i T_{1} \widehat{A}} \tag{6.186}
\end{align*}
$$

for the generic evolution operator of this type. Our aim is to find the timings $T_{i}$ that would result in a desired evolution of the system $\widehat{U}_{d}$.

The identity operator $\widehat{U}_{d}=\widehat{I}$ that corresponds to a non-trivial timing $\sum_{i} T_{i}=T>0$ is the simplest example of a desired evolution. Such an evolution allows one to return the system to its initial state without even knowing this state explicitly. The straightforward way to find this evolution by minimizing the function

$$
\begin{equation*}
\operatorname{Tr}\left[\left(\widehat{U}\left(T_{1}, \ldots, T_{16}\right)-\widehat{I}\right)\left(\widehat{U}\left(T_{1}, \ldots, T_{16}\right)-\widehat{I}\right)^{*}\right] \tag{6.187}
\end{equation*}
$$

in the 16 -dimensional space of the parameters $T_{i}$ in complete analogy with (6.177) requires an already appreciable amount of numerical work. Moreover, this procedure, started from an arbitrary point in the parameter space $\left\{T_{i}\right\}$, usually converges to a local minimum and yields a matrix $\widehat{U}$ which is far from the desired one. In order to be successful, one should start already in a very close vicinity of the solution - for each dimension the value of $T_{i}$ should be within at least a half of the shortest period corresponding to the evolution operators $e^{-i T \widehat{B}}$ or $e^{-i T \widehat{A}}$. For the $N$-dimensional Hilbert space this vicinity has an exponentially small $\sim 2^{-N^{2}}$ volume relative to the total volume $\sim(2 \pi)^{N^{2}}$ of the parameter space. Therefore the functional (6.187) is not a very useful tool, unless we have some other means to gain an insight into the dynamics of the system by narrowing the domain of the solution search.

An alternative algorithm relies on the natural property of complex quantum systems to have a "rigid" spectrum resulting from the interactioninduced repulsion of the energy eigenstates. Let us discuss this in more detail. Consider a time-dependent Hamiltonian $\widehat{H}(t)=\widehat{H}_{0}+\widehat{V}(t)$ that changes periodically in time with a period $T$. As we have already seen in the context of (5.12), the time-dependent operator can be replaced by a time-independent one $\widehat{H}(\tau)-i \frac{\partial}{\partial \tau}$ containing one additional direction $\tau$ in geometric space. Let us perform the Fourier transformation over the variable $\tau$ and make use of the quasienergy representation (6.147) by complete analogy to the quasimomentum representation for a quantum system, periodic in space. The Hamiltonian adopts the form $\sum_{k} \widehat{H}(k \omega) \delta_{n+k, n^{\prime}}+n \omega \delta_{n, n^{\prime}}$ where $n$ and $n^{\prime}$ are the quasienergy indices and $\omega=2 \pi / T$. For the operator of the evolution over a period
we therefore obtain $\widehat{U}(T)=\exp \left[-i T \sum_{k} \widehat{H}(k \omega) \delta_{n+k, n^{\prime}}\right]$ whereas the term $-i T n \omega \delta_{n, n^{\prime}}$ in the Hamiltonian just gives a multiple of $2 \pi$ phase for each $n$ and can be omitted. This implies that the energy eigenstates of the effective Hamiltonian $i T^{-1} \ln \widehat{U}(T)$ separated by the energy $\omega$ of a quantum or by a multiple of this value, have to be considered as neighboring states from the point of view of the contributions they give to the operator of evolution over a period. In other words, the eigenstates of $\widehat{H}_{0}$ located at an energy distance $n \omega$ are coupled by the $n$-th harmonics of the time-dependent part $\widehat{V}(t)$ and repel each other in the same way as the stationary coupled states that lie close together in energy. Mutual repulsion of the levels causes the eigenvalues of $\widehat{U}(T)$ to be distributed over the unit circle in the complex plane with the maximum possible distance among the neighboring eigenvalues. We illustrate this in Fig. 6.20(b).

In the context of two interacting, two-level systems we now consider the truncated evolution operator

$$
\begin{equation*}
\widehat{U}_{0}\left(\left\{T_{i}\right\}\right)=e^{-i T_{4} \widehat{B}} e^{-i T_{3} \widehat{A}} e^{-i T_{2} \widehat{B}} e^{-i T_{1} \widehat{A}} \tag{6.188}
\end{equation*}
$$

which depends on $N=4$ parameters. All four eigenvalues of this $4 \times 4$ matrix are on the unit circle in the complex plane. Let us try to choose the four exposure times $\left\{T_{i}\right\}$ such that the eigenvalues of $\widehat{U}_{0}$ become equally spaced on the unit circle, thus aiding this natural tendency, typical of strong and complex perturbations. For this purpose we calculate the coefficients $c_{n}\left(\left\{T_{i}\right\}\right)$ of the characteristic polynomial

$$
\begin{align*}
\operatorname{det}\left|\widehat{U}_{0}\left(\left\{T_{i}\right\}\right)-\lambda \widehat{I}\right|= & \lambda^{4}+\lambda^{3} c_{3}\left(\left\{T_{i}\right\}\right)+\lambda^{2} c_{2}\left(\left\{T_{i}\right\}\right) \\
& +\lambda c_{1}\left(\left\{T_{i}\right\}\right)+c_{0}\left(\left\{T_{i}\right\}\right) \tag{6.189}
\end{align*}
$$

and minimize the sum

$$
\begin{equation*}
F\left(\left\{T_{i}\right\}\right)=\sum_{n=1}^{n=N-1}\left|c_{n}\left(\left\{T_{i}\right\}\right)\right|^{2} \tag{6.190}
\end{equation*}
$$

of its coefficients. When the function (6.188) assumes a zero value, the roots of the resulting characteristic polynomial $\lambda^{N}+c_{0}=0$ are equidistantly distributed on the unit circle. Due to the fact that a homogeneous distribution of the eigenvalues of $\widehat{U}_{0}$ is a natural property of complex systems, the function (6.190) reaches the global minimum much more often, as compared to the function (6.187).

Being repeated $N=4$ times, the truncated evolution yields the identity transformation $\widehat{U}_{0}^{4}=\widehat{I}$. Starting from this non trivial identity $\widehat{U}\left(\left\{T_{i}\right\}\right)$ with $\left\{T_{i}\right\}=T_{1, \ldots,}, T_{4}, T_{1}, \ldots, T_{4}, T_{1}, \ldots, T_{4}, T_{1}, \ldots, T_{4}$ one finds the desired transformation $\widehat{U}_{d}$ by solving iteratively (with the help of $(6.181),(6.182)$ ) the problem $\widehat{U}\left(\left\{T_{i}\right\}\right)=\exp \left(\mu \ln \widehat{U}_{d}\right)$ for $\mu$ gradually increasing from 0 to 1 .

### 6.6.3 Control of Coherence Loss

Decay and loss of coherence are two processes, common to all open quantum systems interacting with an environment. Two different relaxation times $T_{1}$ and $T_{2}$ characterize these processes in two-level systems as has been mentioned in (5.60). Though both processes represent irreversible relaxation, they are of rather different nature. The number of occupied quantum states, given by the rank of the density matrix - an analog of the phase volume in the classical case - may change in an elementary act of relaxation, and usually does it towards diminishing. By contrast, this number does not change in an elementary act of coherence loss, but yet increases after being averaged over subsequent elementary acts.

Considering irreversible processes such as decay or measurement, one can always employ certain physical models, describing these phenomena as a coherent process in a large compound system, that apart from the system of interest, also includes a much larger part with a continuous spectrum called "the environment" and "the meter", respectively.- The most relevant model for the measurement and for the decay is the level-continuum problem considered on p. 72. Here we will not dwell on this aspect of the process, simply assuming that after a given time considered as short relative to other typical times, the system of interest simply moves to a certain state (either known of unknown) with an exponentially vanishing probability of the opposite result.

On the contrary, loss of coherence does not require a quantum model of the environment. It can be viewed as a systematic and uncontrolled classical external intervention in the system, resulting in the application of a known Hamiltonian during unknown time intervals. Without performing a measurement, one can never know whether such an intervention has occurred. Coherence loss protection is therefore a special kind of measurement which produces the action of canceling all of the results of coherence loss without affecting the quantum state of the system.

We now show how one can completely suppress coherence loss in a universal way. The suggested method relies on the basic ideas of classical errorcorrecting codes and on the idea known as the Zeno effect, that states that periodic measurements allow one to keep a quantum system in the initial state for as long as needed, provided that the period $T$ is short relative to the relaxation time $T_{2}$. In order to protect a quantum system from coherence loss, we do not intervene in the entire system, but only affect the auxiliary part, the ancilla, which we have to add for the purpose of protection. We perform encoding in the entire Hilbert space of the system and the ancilla by employing non-holonom control.

We start with a geometric illustration of the regular Zeno effect for the simplest two-level system, show how this method allows one to inhibit the coherence loss, and suggest a generalization to larger systems. The density operator of a two-level system is a $2 \times 2$ matrix which can be represented according to (6.166) as a sum of four parts $\widehat{\rho}=\widehat{I}+P \widehat{\sigma}_{x}+Q \widehat{\sigma}_{y}+W \widehat{\sigma}_{z}$ where
three Pauli matrices are multiplied by three real components $P, Q$, and $W$ of the Bloch vector. For pure states, the rank of $\widehat{\rho}$ is 1 and the Bloch vector remains on the unit sphere. When the density matrix becomes a statistical mixture of different pure states, the vector becomes shorter. The coherence loss, associated with the action of the unknown Hamiltonian, leaves the vector on the unit sphere, although moving it to a new, unknown position as shown in Fig. 6.21(a). Usually the coherence loss time $T_{2}$ is much shorter than the


Fig. 6.21. The Zeno effect for a two-level system. (a) The Bloch vector $B$ representing the density matrix of a two-level system in terms of the polarization $P$, dispersion $Q$, and the population inversion $W$. Pure states correspond to the unit sphere. The Hamiltonian evolution, including the uncontrolled loss of coherence with a time constant $T_{2}$, occurs on this surface. Decay destroys pure states, and moves the Bloch vector out of the surface with a typical time $T_{1}$ (dash-dot line). (b) The trajectory of the Bloch vector in the $P-W$ plane. According to first-order time-dependent perturbation theory the coherence loss (dashed arrow) and decay to a narrow band (solid curve) do not move the Bloch vector away from the unit sphere, and can therefore be corrected to second order by frequent measurementinduced projections (dotted arrow) onto the initial upper state $W=1$. The decay (dash-dot line) moves $B$ out of the sphere with a typical time $T_{1}$.
longitudinal relaxation time $T_{1}$ which is associated with spontaneous decay and entanglement, with an escaping photon thus moving the Bloch vector out of the unit sphere. In the case of spontaneous decay to an infinite continuum, this process results in a linear in time decrease of the length of the Bloch vector. For a band-like continuum, a short period of quadratic decrease of the length may exist prior to this linear regime. This persists for a time equal to the inverse width of the band. These three possibilities are shown in Fig. 6.21(b).

Measurement of a physical quantity results in the projection of the Bloch vector to a vector representing the measured eigenstate of the corresponding
operator. In Fig. 6.21 we illustrate this for the measurement of the population of the upper level, corresponding to the positive eigenvalue of the operator $\widehat{\sigma}_{z}$. Being subjected to the action of a Hamiltonian of type $\widehat{\sigma}_{y}$, for unknown time $T$, the Bloch vector turns from its initial upper position in the direction $P$ to an unknown angle $\sim T / T_{2}$ which is small, provided that $T \ll T_{2}$. The linear increment of the vector, resulting from the evolution, is therefore orthogonal to the vector itself. Being projected after the measurement to the $W$-axis, it returns to the initial position almost entirely. In other words, the regular Zeno effect allows one to restore the initial state vector up to the second-order terms in $T / T_{2}$. Though for a multilevel system, the Bloch representation is less convenient, the essence of the Zeno effect relies on the same geometric idea. We consider the real and the imaginary parts of the probability amplitudes as independent components of a real state vector $v$, which rotates in the course of an uncontrolled unitary evolution at a small angle, and returns after the measurement to the initial position, up to the second-order corrections.

In the $N$-dimensional Hilbert space $\mathcal{H}_{\mathcal{N}}$ of a compound quantum system $\mathcal{N}$ comprising the main part $\mathcal{K}$ and an axillary part, the ancilla $\mathcal{A}=\mathcal{N} / \mathcal{K}$, a measurement of the ancilla state results in the projection of the compound state vector onto a multidimensional subspace $\mathcal{H}_{\mathcal{K}}$ instead of a state vector. The dimension $K$ of this subspace coincides with the dimension of the Hilbert space of the main part. If the evolution induced by the error Hamiltonians results in a linear increment of the state vector, always perpendicular to $\mathcal{H}_{\mathcal{K}}$, then the projecting measurement returns the system to its initial state. This is sketched in Fig. 6.22(a) in complete analogy to the regular Zeno effect. In other words, the errors do not affect the main part of the system. Though this assumption corresponds to an unphysical situation, when all errors affect only the ancilla and no error Hamiltonians act on the main system, it nevertheless allows us to make an important conclusion that the total number $M$ of different error-inducing Hamiltonians may not exceed the dimensionality $A$ of the ancilla Hilbert space. Universal protection against this number of Hamiltonians relies on a similar method, which just requires some modification in the spirit of coding theory.

We come to the key idea of the protection method. For a system composed of a main part and an ancilla, eventually subjected to an action of a number $M$ of different but known error-inducing Hamiltonians $E_{m}$, one can find a subspace $\mathcal{C}$ such that any error-induced linear increment $\delta v$ of any state vector $v \in \mathcal{C}$ is orthogonal to $\mathcal{C}$, provided $M \leq A$. Consider the subspace $\mathcal{H}_{\mathcal{K}}$, formed as the direct product of the Hilbert space of the main part by a given state vector $|\widetilde{\alpha}\rangle$ of the ancilla. If by a unitary coding transformation $\widehat{C}$, one can rapidly move $\mathcal{H}_{\mathcal{K}}$ to $\mathcal{C}$ prior to exposing the system to the action of errors, and perform rapidly the decoding transformation $\widehat{C}^{-1}$ just before measurement of the ancilla in the state $|\widetilde{\alpha}\rangle$, then the resulting dynamics are
equivalent to the errors acting solely on the ancilla. Hereafter we call $\mathcal{C}$ a code space. A system $\mathcal{N}$ may have many $\mathcal{C}$.


Fig. 6.22. (a) The multidimensional Zeno effect. If all the linear increments $\delta v \sim \widehat{E}_{j} v$ of the state vectors $v$ from the code subspace $\mathcal{C}$ are orthogonal to this subspace for any error-induced transformation $v^{\prime} \rightarrow v$, the projection to $\mathcal{C}$ induced by the measurement compensates for the errors up to $O\left(\left(T / T_{2}\right)^{2}\right)$. (b) The iterative algorithm for finding the code space. The code space $\mathcal{C}$ is represented here by a single vector. The application of all possible errors $\widehat{E}_{m} \mathcal{C}^{\prime}$ to a code space guess $\mathcal{C}^{\prime}$ forms the space of errors $\mathcal{E}^{\prime}$. A linear combination $\sum \gamma_{m} \widehat{E}_{m} \mathcal{C}^{\prime}$ shows a direction in $\mathcal{E}^{\prime}$ in which the length of the sum $\mathcal{C}^{\prime}+\sum \gamma_{m} \widehat{E}_{m} \mathcal{C}^{\prime}$ assumes the minimum value. Displacement in this direction increases the angle $\Theta$ between $\mathcal{C}^{\prime}$ and $\mathcal{E}^{\prime}$, which approaches the maximum value $\pi / 2$ in the vicinity of the "median" point. Repeating this exponentially converging iteration one obtains the code space $\mathcal{C}$ orthogonal to the error space $\mathcal{E}$.

Let us perform formal calculations illustrating this idea. Consider a state vector $|S\rangle=|s\rangle \otimes|\widetilde{\alpha}\rangle$ of $\mathcal{N}$ formed by $\mathcal{K}$ in a state $|s\rangle$ and $\mathcal{A}$ in the state $|\widetilde{\alpha}\rangle$. We treat the errors $\widehat{E}_{m}$ as matrix parts of interaction Hamiltonians with external random fields $f_{m}(t)$ that produce an uncontrolled evolution of the system. For $\widehat{E}_{m} \sim 1$ one can estimate $f_{m} \sim 1 / T_{2}$. Under the action of the error-inducing Hamiltonians $\widehat{E}_{m} f_{m}(t)$, the vector $|S\rangle$ undergoes an uncontrolled evolution given by the unitary operator

$$
\begin{equation*}
\widehat{U}=\prod_{m=1}^{M} e^{-i \widehat{E}_{m} \int f_{m}(t) d t} \simeq \widehat{I}-i \sum_{m=1}^{M} \widehat{E}_{m} \int f_{m}(t) d t \tag{6.191}
\end{equation*}
$$

where the time ordering of the product is implicit. Actions over the Zeno period $T$ of the $M$ different fields are small $\int_{t}^{t+T} f_{m}(x) d x\left|\widehat{E}_{m}\right| \ll 1$, such that only the identity operator $\widehat{I}$ and the first-order terms $\sim \int f_{m}(t) d t$ are important in the Taylor series for the evolution operator.

Let $\widehat{\rho}_{s c}=|s\rangle \otimes\langle s|$ be the density matrix of the main part before the action of errors. Then $\widehat{\rho}=|S\rangle \otimes\langle S|=|\widetilde{\alpha}\rangle \widehat{\rho}_{s c}\langle\widetilde{\alpha}|$ is the density matrix of the entire system. After coding by the unitary transformation $\widehat{C}$, this matrix takes
the form $\widehat{C} \widehat{\rho} \widehat{C}^{-1}$, and reads $\widehat{U} \widehat{C} \widehat{\rho} \widehat{C}^{-1} \widehat{U}^{*}$ at the end of the action of errors. Decoding yields $\widehat{\rho}^{\prime}=\widehat{C}^{-1} \widehat{U} \widehat{C} \widehat{\rho} \widehat{C}^{-1} \widehat{U}^{*} \widehat{C}$ whereas the direct substitution shows that the variation $\delta \widehat{\rho}_{s c}=\langle\widetilde{\alpha}| \widehat{\rho}^{\prime}-\widehat{\rho}|\widetilde{\alpha}\rangle$ of the density matrix of the main system after the perturbation (6.191) and the measurement of the ancilla is given by the commutator

$$
\begin{equation*}
\delta \widehat{\rho}_{s c}=-i\left[\sum_{m=1}^{M} \int f_{m}(t) d t\langle\widetilde{\alpha}| \widehat{C}^{-1} \widehat{E}_{m} \widehat{C}|\widetilde{\alpha}\rangle, \widehat{\rho}_{s c}\right] . \tag{6.192}
\end{equation*}
$$

Therefore $\widehat{\rho}_{s c}$ satisfies the equation

$$
\begin{equation*}
i \frac{d \widehat{\rho}_{s c}}{d t}=\left[\widehat{h}_{e}, \widehat{\rho}_{s c}\right] ; \widehat{h}_{e}=\sum_{m=1}^{M} f_{m}\langle\widetilde{\alpha}| \widehat{C}^{-1} \widehat{E}_{m} \widehat{C}|\widetilde{\alpha}\rangle \tag{6.193}
\end{equation*}
$$

with the effective Hamiltonian $\widehat{h}_{e}$ made up of error-inducing Hamiltonians transformed by coding and decoding and projected to the state of the ancilla. Our aim is to set these projections to zero by a proper choice of the coding transformation $\widehat{C}$.

The relation $|v\rangle=\widehat{C}|s\rangle \otimes|\widetilde{\alpha}\rangle$ holds for the code vectors and the states of the system. If $\widehat{h}_{e}=0$, the errors do not affect the main part of the system and $\widehat{\rho}_{s c}=$ const which implies

$$
\begin{equation*}
\left\langle v^{\prime}\right| \widehat{E}_{m}|v\rangle=0 ; \quad \forall|v\rangle,\left|v^{\prime}\right\rangle \in \mathcal{C} ; \quad \forall \widehat{E}_{m} \in\left\{\widehat{E}_{m}\right\} \tag{6.194}
\end{equation*}
$$

In other words, a linear increment $\delta v=\widehat{E}_{m} v$ resulting from any error $\widehat{E}_{m}$ applied to a vector $v$ of the coding space is orthogonal to any other vector $v^{\prime}$ of the coding space. Note that all Hamiltonians $\widehat{E}_{m}$ in (6.194) projected to the code space have zero traces, or should be set in this form by subtracting a scalar operator $\lambda \widehat{I}$ that does not change the Hamiltonian dynamics. We note the physical meaning of (6.194). A corrupted state vector of the code space does not overlap with another vector of the code space if they were orthogonal before the corruption. This corresponds to the so-called Hamming error detection condition in classical information theory. It simply means that the variety of all of the results obtained by the corruption of one message should not include another possible message. The only difference is that in quantum mechanics we interpret the overlap as a scalar product of the state vectors in the Hilbert space, whereas in classical information theory this concept is formulated in terms of "distance", a metrical property which denotes the number of distinct information bits for two messages.

Now we show an algorithm which allows one to design such codes. For $N=K A$ dimensional Hilbert space $\mathcal{H}_{\mathcal{N}}$ formed by the main part $\mathcal{K}$ and ancilla $\mathcal{A}$, condition (6.194) is a set of $M K^{2}$ equations taken over $A K^{2}$ matrix elements $\langle\alpha| \otimes\left\langle s^{\prime}\right| \widehat{C}|s\rangle \otimes|\widetilde{\alpha}\rangle$ of the $N \times N$ encoding operator $\widehat{C}$. This defines a rectangular $K \times N$ part $\widehat{C}|\widetilde{\alpha}\rangle$ of $\widehat{C}$. To specify (6.194) further, let the states
$|s\rangle \otimes|\widetilde{\alpha}\rangle$ correspond to the first $K$ positions of the state vector, which implies that each of the $M$ matrices $\widehat{C}^{-1} \widehat{E}_{m} \widehat{C}$ has an all-zero upper left $K \times K$ corner. The standard methods of linear algebra do not give a recipe for finding a linear transformation $\widehat{C}$ that simultaneously sets to zero the corners of $M$ different matrices. The following iterative algorithm accomplishes this task- we find a code space $\mathcal{C}$ composed of all linear combinations of $K$ first columns of $\widehat{C}$ which satisfies the conditions of (6.194) given $M$ arbitrary error Hamiltonians $\widehat{E}_{m}$.

Though it can directly be done in the Hilbert space, in order to elucidate the geometrical meaning of the algorithm we write the code basis as a single real vector C of length $2 K N$ by appending vectors of length $N$ corresponding to real and imaginary parts of the first $K$ columns of $\widehat{C}$. We rewrite (6.194) in the form

$$
\begin{equation*}
\left(\mathrm{C}_{m} \mathrm{C}\right)=0 \tag{6.195}
\end{equation*}
$$

where $2 K N \times 2 K N$ error matrices $\widehat{\mathrm{E}}_{m}$ have a $K \times K$ block structure. Each diagonal and upper-diagonal block repeats real $2 N \times 2 N$ matrices $\widehat{\mathrm{E}}_{m}$ corresponding to Hamiltonians $\widehat{E}_{m}$, whereas the lower blocks are formed by the unity operator $\widehat{I}$ in the Hilbert space $\mathcal{H}_{\mathcal{N}}$ rewritten in the chosen real representation. The first type of blocks corresponds to the condition (6.194) whereas the second part ensures the orthogonality of different columns of $\widehat{C}|\widetilde{\alpha}\rangle$. In other words, the vector $C$ should satisfy $M$ quadratic equations (6.195) for each $m$.

The input of our algorithm is a randomly selected and properly normalized vector $C^{\prime}$. If by chance $C^{\prime}$ is orthogonal to all vectors $\widehat{\mathrm{E}}_{m} \mathrm{C}^{\prime}$, the problem is solved. Otherwise, there exists a linear combination $\sum_{m=1}^{M} \gamma_{m} \widehat{\mathrm{E}}_{m} \mathrm{C}^{\prime}$ that minimizes the length of the vector $\sum_{m=1}^{M} \gamma_{m} \widehat{\mathrm{E}}_{m} \mathrm{C}^{\prime}+\mathrm{C}^{\prime}$. This is schematically shown in Fig.6.22(b). To find it, we employ the standard variational principle and solve the corresponding set of linear equations for $\gamma_{m}$. When we displace the end of our vector $C^{\prime}$ in this direction, the angle $\Theta$ between this vector and the plane $\left\{\widehat{\mathrm{E}}_{m} \mathrm{C}^{\prime}\right\}$ increases. Since the orientation of the plane $\left\{\widehat{\mathrm{E}}_{m} \mathrm{C}^{\prime}\right\}$ may also change with the displacement, the latter should not be too big. To achieve the best convergency of our further iterations we direct the new vector to the "median" point $\mathrm{C}^{\prime} \rightarrow \sum_{m=1}^{M} \frac{1}{2} \gamma_{m} \widehat{\mathrm{E}}_{m} \mathrm{C}^{\prime}+\mathrm{C}^{\prime}$, then normalize the new vector, and repeat this procedure until the required accuracy is achieved.

Summarizing, we state that the discussed protecting codes give a universal means to protect a quantum system from coherence loss, although they require a physical procedure of rapid coding and frequent measurement which determines the efficiency of the error protection.

## 7 The Dynamics of One-Dimensional Relay-Type Systems

Thus far we have mainly considered an extreme of quantum complexity, where each pair of quantum states has been coupled by a matrix element, which is statistically equivalent to any matrix element coupling another pair of levels. Although this was not strictly speaking the case for two-band systems where states of the same bands are not coupled, the resulting impact of this fact on their behavior is so tiny that we could attribute them to the same class of problems. The structure of the couplings among different quantum states is closely related to the spatial dimensionality of the system, and the situation where all of the coupling matrix elements have the same order of magnitude is often referred to as the 0-dimensional case. In this chapter we consider a different situation, where quantum systems are arranged, satisfying certain "selection rules". We begin with the consideration of multilevel systems where not all of the states are coupled to each other, but where this coupling is organized in some way. The first natural example is a system of isolated levels, coupled only with their closest neighbors. After presenting several examples of exactly soluble problems, we consider systems conforming to the approximative WKB description, and then turn to the systems with randomly distributed couplings and energy level positions. This will give us an opportunity to introduce a number of mathematical methods useful for the analysis of complex systems.

### 7.1 Exactly Soluble Relays of Isolated Levels

Consider a relay of isolated levels, that is a set of quantum states $|n\rangle$ of energies $E_{n}$ coupled with their closest neighbors by the interaction matrix elements $V_{n, m}=V_{n} \delta_{n, m+1}+V_{n}^{*} \delta_{m+1, n}$, where $\delta_{m, n}$ is the Kronecker delta. The Schrödinger equation for the probability amplitudes $\psi_{n}$ corresponding to the states $|n\rangle$ reads

$$
\begin{equation*}
i \frac{\partial \psi_{n}}{\partial t}=E_{n} \psi_{n}+V_{n} \psi_{n-1}+V_{n+1}^{*} \psi_{n+1} \tag{7.1}
\end{equation*}
$$

The corresponding equation for the eigenvalues

$$
\begin{equation*}
0=\Delta_{n} \psi_{n}+V_{n} \psi_{n-1}+V^{*}{ }_{n+1} \psi_{n+1} \tag{7.2}
\end{equation*}
$$

with $\Delta_{n}=E_{n}-E$ has the form of a three-term recurrence relation, typical of orthogonal polynomials or some other special functions. Each of these special functions corresponds to an exactly soluble problem for the particular case of $\Delta_{n}$ and $V_{n}$. We consider a number of such problems.

### 7.1.1 Uniform Coupling and Linear Detuning

The simplest case $\Delta_{n}=\alpha n-E, V_{n}=V_{n}^{*}=V$ conforms to the recurrence relations for the Bessel functions

$$
\begin{equation*}
0=(\alpha n-E) J_{n-\frac{E}{\alpha}}\left(\frac{-V}{2 \alpha}\right)+V J_{n-1-\frac{E}{\alpha}}\left(\frac{-V}{2 \alpha}\right)+V J_{n+1-\frac{E}{\alpha}}\left(\frac{-V}{2 \alpha}\right) \tag{7.3}
\end{equation*}
$$

and therefore the energy eigenstates can indeed be given in terms of the Bessel functions

$$
\begin{equation*}
\psi_{n}(E)=A J_{n-\frac{E}{\alpha}}\left(\frac{-V}{2 \alpha}\right)+B J_{-n+\frac{E}{\alpha}}\left(\frac{V}{2 \alpha}\right) \tag{7.4}
\end{equation*}
$$

where $A$ and $B$ are constants.
Dealing with the time-dependent problem for the initial condition $\psi_{0}(t=$ $0)=1$, it is however easier to make use of the generating function $\Psi(t, \varphi)=$ $\sum_{n} \psi_{n}(t) e^{i \varphi n}$ for which, for the chosen dependencies of $\Delta_{n}$ and $V_{n}$ after substitution into (7.1), one finds the equation

$$
\begin{align*}
i \frac{\partial \Psi(t, \varphi)}{\partial t} & =\sum_{n} \alpha n \psi_{n}(t) e^{i \varphi n}+\left(V e^{i \varphi}+V^{*} e^{-i \varphi}\right) \Psi(t, \varphi) \\
& =-i \alpha \frac{\partial \Psi(t, \varphi)}{\partial \varphi}+\left(V e^{i \varphi}+V^{*} e^{-i \varphi}\right) \Psi(t, \varphi) \tag{7.5}
\end{align*}
$$

which can also be extended to the case of a time-dependent coupling $V(t)$ and results in

$$
\begin{equation*}
i \frac{\partial \ln [\Psi(t, \varphi)]}{\partial t}+i \alpha \frac{\partial \ln [\Psi(t, \varphi)]}{\partial \varphi}=V(t) e^{i \varphi}+V^{*}(t) e^{-i \varphi} \tag{7.6}
\end{equation*}
$$

We now introduce the new variables $x=(\alpha t+\varphi) / 2$ and $y=(\alpha t-\varphi) / 2$ which yields

$$
\begin{equation*}
\ln [\Psi(t, \varphi)]=\int V\left(\frac{x+y}{\alpha}\right) e^{i(x-y)}+V^{*}\left(\frac{x+y}{\alpha}\right) e^{i(y-x)} \frac{d x}{i \alpha} \tag{7.7}
\end{equation*}
$$

With the allowance for the initial condition $\Psi(t=0, \varphi)=1$ one therefore finds

$$
\begin{align*}
\ln [\Psi(t, \varphi)] & =\int_{\varphi / 2}^{(\alpha t+\varphi) / 2} V\left(\frac{x+y}{\alpha}\right) e^{i(x-y)}+V^{*}\left(\frac{x+y}{\alpha}\right) e^{i(y-x)} \frac{d x}{i \alpha} \\
& =\int_{t / 2}^{t} V(z) e^{i(\alpha z-\alpha t+\varphi)}+V^{*}(z) e^{-i(\alpha z-\alpha t+\varphi) \frac{d z}{i}}, \tag{7.8}
\end{align*}
$$

where the replacement $z \rightarrow(x+y) / \alpha$ has been made.

From the generating function one immediately finds the amplitudes

$$
\begin{equation*}
\psi_{n}(t)=\int_{0}^{2 \pi} \frac{d \varphi}{2 \pi} \exp \left\{-i\left[e^{i \varphi} \int_{t / 2}^{t} V(z) e^{i(\alpha z-\alpha t)} d z+c . c .\right]-i n \varphi\right\} \tag{7.9}
\end{equation*}
$$

which, with the help of the integral representation for the Bessel functions

$$
\begin{equation*}
J_{n}(X)=\frac{1}{2 \pi} \int_{-\pi}^{\pi} e^{-i X \sin \vartheta+i n \vartheta} d \vartheta \tag{7.10}
\end{equation*}
$$

yields

$$
\begin{equation*}
\psi_{n}(t)=i^{n} e^{i n \arg Z} J_{n}(2|Z|) \tag{7.11}
\end{equation*}
$$

where $Z=\int_{t / 2}^{t} V(z) e^{i(\alpha z-\alpha t)} d z$.
We note that the problem can also be solved for any dependence $\alpha(t)$ which results in (7.11) with $Z=\int_{t / 2}^{t} V(z) \exp \left[i \int_{t}^{z} \alpha(\theta) d \theta\right] d z$. The reason why this can be solved for arbitrary dependences $V(t)$ and $\alpha(t)$ lies in the algebraic structure of the Hamiltonian, which only contains the operators $\widehat{a}_{ \pm} \rightarrow\left(\delta_{n, m+1} \pm \delta_{m+1, n}\right) / 2$ and $\widehat{a}_{d} \rightarrow n \delta_{n, m}$ forming a closed subalgebra in the infinite dimensional Hilbert space of the entire system. These operators satisfy the commutation relations $\left[\widehat{a}_{+}, \widehat{a}_{-}\right]=0,\left[\widehat{a}_{d}, \widehat{a}_{ \pm}\right]=\widehat{a}_{\mp}$ which is easy to see after Fourier transformation over the indices (momentum representation $n \rightarrow \varphi$ ), since $\widehat{a}_{d} \rightarrow-i \partial / \partial \varphi ; \widehat{a}_{+} \rightarrow \cos \varphi$, and $\widehat{a}_{-} \rightarrow i \sin \varphi$. That is why the representation (5.5) for the evolution operator of this infinite-dimensional system contains only these three generators, demonstrating once again the fact that an exactly soluble problem usually has a hidden operator structure symmetry.

In Fig. 7.1(a) we show the population distribution corresponding to the amplitudes (7.11) for different values of the parameter $Z$. Note that the states with $n<V$, apart from periodic oscillations, are populated almost uniformly, with a certain augmentation for $n \simeq V$, whereas for $n>V$ the population decreases exponentially.

### 7.1.2 The Harmonic Oscillator in an Arbitrary Time-Dependent Field

The next important case of (7.1) is $E_{n}=\alpha n, V_{n}=V \sqrt{n}$ which after the replacement $\psi_{n}=e^{-i \alpha n t} \widetilde{\psi}_{n}$ results in

$$
\begin{equation*}
i \frac{\partial}{\partial t} \widetilde{\psi}_{n}=V e^{i \alpha t} \sqrt{n} \widetilde{\psi}_{n-1}+V^{*} e^{-i \alpha t} \sqrt{n+1} \widetilde{\psi}_{n+1} \tag{7.12}
\end{equation*}
$$

which it is expedient to rewrite with the help of the operators $\widehat{a}$ and $\widehat{a} \dagger$ of the harmonic oscillator


Fig. 7.1. Population distribution in systems with closest neighbor interaction. (a) Linear detuning and constant coupling. The absolute value of the amplitude is given by the Bessel functions (7.11). (b) Harmonic oscillator. The absolute value of the amplitude is given by the Laguerre polynomials (7.20). The dependence is shown for the initially populated state $n^{\prime}=6$.

$$
\begin{equation*}
i \frac{\partial}{\partial t}|\widetilde{\psi}\rangle=V e^{i \alpha t} \widehat{a}|\widetilde{\psi}\rangle+V^{*} e^{-i \alpha t} \widehat{a} \dagger|\widetilde{\psi}\rangle \tag{7.13}
\end{equation*}
$$

This form of the equation will help us to find a solution even for an arbitrary time dependence of $V(t)$ and $\alpha(t)$ when (7.13) reads

$$
\begin{equation*}
i \frac{\partial}{\partial t}|\widetilde{\psi}\rangle=\left[V(t) e^{i \int \alpha(t) d t} \widehat{a}+V^{*}(t) e^{-i \int \alpha(t) d t} \widehat{a}^{\dagger}\right]|\widetilde{\psi}\rangle \tag{7.14}
\end{equation*}
$$

We introduce the quantity

$$
\begin{equation*}
Z(t)=-i \int V(t) e^{i \int \alpha(t) d t} d t \tag{7.15}
\end{equation*}
$$

and consider a transformation $\widehat{U}(t)=\exp \left[Z(t) \widehat{a}+Z^{*}(t) \widehat{a} \dagger\right]$, which according to the Baker-Campbell-Hausdorff formula (5.3) and with allowance of the commutation relation $[\widehat{a}, \widehat{a} \dagger]=1$ can be written in the form

$$
\begin{align*}
\widehat{U}(t) & =\exp \left[Z(t) \widehat{a}+Z^{*}(t) \widehat{a} \dagger\right] \\
& =\exp [Z(t) \widehat{a}] \exp \left[Z^{*}(t) \widehat{a} \dagger\right] \exp \left[\frac{1}{2} Z(t) Z^{*}(t)\right] . \tag{7.16}
\end{align*}
$$

For $|\widetilde{\psi}\rangle=\widehat{U}(t)\left|\psi^{\prime}\right\rangle$ after straightforward calculations, (7.14) yields

$$
\begin{equation*}
\frac{\partial}{\partial t}\left|\psi^{\prime}\right\rangle=\left[Z(t) \frac{\partial Z^{*}(t)}{2 \partial t}-Z^{*}(t) \frac{\partial Z(t)}{2 \partial t}\right]\left|\psi^{\prime}\right\rangle \tag{7.17}
\end{equation*}
$$

where in the course of calculations the expansion

$$
\begin{equation*}
e^{\widehat{X}} \widehat{Y} e^{-\widehat{X}}=\widehat{Y}+[\widehat{X}, \widehat{Y}]+\frac{1}{2!}[\widehat{X},[\widehat{X}, \widehat{Y}]]+\ldots \tag{7.18}
\end{equation*}
$$

has been employed.

Equation (7.17) describes an unimportant common phase shift, and hence the operator $\widehat{U}(t)$ of (7.16), depending on two real functions $\operatorname{Re} Z(t)$ and $\operatorname{Im} Z(t)$, gives the evolution of the system. However, if the initial state of the system corresponds to a single number state, that is $\psi_{n}(t=0)=\delta_{n, n^{\prime}}$, and if we are interested only in the population distribution and not in the relative phases of different probability amplitudes, the result depends only on a single function $|Z(t)|$ and reads

$$
\begin{equation*}
\left.\left|\psi_{n}\right|^{2}=\left|\langle n| \exp \left[|Z(t)|\left(\widehat{a}-\widehat{a}^{\dagger}\right)\right]\right| n^{\prime}\right\rangle\left.\right|^{2} \tag{7.19}
\end{equation*}
$$

which can be obtained by the transformation

$$
\exp \left(-i \Phi(t) \widehat{a}^{\dagger} \widehat{a}\right) \widehat{U}(t) \exp \left(i \Phi(t) \widehat{a}^{\dagger} \widehat{a}\right)
$$

corresponding to rotation of the operator (7.16) by an angle $\Phi(t)=\pi / 2+$ $\arg Z(t)$ in the phase space. In the combination $\left(\widehat{a}-\widehat{a}^{\dagger}\right) / \sqrt{2}$ one recognizes the momentum operator $i \widehat{p}=\partial / \partial x$, and therefore $\exp \left[|Z(t)|\left(\widehat{a}-\widehat{a}^{\dagger}\right)\right]$ is just an operator of finite shift (5.13) of the coordinate by the value $|Z(t)| \sqrt{2}$. The simplest way to determine the matrix elements (7.19) of this operator is simply to perform calculations in the coordinate representation, where the wave functions on the oscillator energy eigenstates $\varphi(x)$ are given by Hermite polynomials $H_{n}(x) e^{-x^{2} / 2} / \sqrt{2^{n} n!}$. This yields

$$
\begin{align*}
\left|\psi_{n}\right|^{2} & =\frac{e^{-|Z(t)|^{2}}}{2^{n+n^{\prime}} n!n^{\prime}!}\left|\int H_{n^{\prime}}\left(x+\frac{|Z(t)|}{\sqrt{2}}\right) H_{n}\left(x-\frac{|Z(t)|}{\sqrt{2}}\right) e^{-x^{2}} d x\right|^{2} \\
& =\pi e^{-|Z(t)|^{2}}|Z(t)|^{2\left(n-n^{\prime}\right)} \frac{n^{\prime}!}{n!}\left[L_{n^{\prime}}^{n-n^{\prime}}\left(|Z(t)|^{2}\right)\right]^{2} \tag{7.20}
\end{align*}
$$

where $n>n^{\prime}$ and $L_{k}^{m}(x)$ is a Laguerre polynomial. For $n>n^{\prime}$ these two indices have to be interchanged.

In Fig. 7.1(b) we show the population distribution corresponding to the amplitudes (7.20) for different values of the parameter $Z$. The existence of an explicit solution again originates in the structure of the operator algebra corresponding to the problem. The Hamiltonian for the harmonic oscillator, interacting with an external field, can be given in terms of three operators - the unperturbed Hamiltonian $\widehat{H}_{0}=\widehat{a}^{\dagger} \widehat{a}=\left(\widehat{P}^{2}+\widehat{Q}^{2}\right) / 2$, the momentum $\widehat{P}=i\left(\widehat{a}^{\dagger}-\widehat{a}\right) / \sqrt{2}$ and the coordinate $\widehat{Q}=\left(\widehat{a}+\widehat{a}^{\dagger}\right) / \sqrt{2}$ operators, which satisfy the canonical commutation relations $\left[\widehat{H}_{0} \widehat{Q}\right]=i \widehat{P},\left[\widehat{H}_{0} \widehat{P}\right]=-i \widehat{Q}$, and $[\widehat{P} \widehat{Q}]=-i$.

### 7.1.3 Raman Pumping of a Harmonic Oscillator

The case $E_{n}=2 \alpha n, V_{n}=V \sqrt{2 n(2 n+1)}$, where $V$ and $\alpha$ are functions of time, also corresponds to a harmonic oscillator, although subjected to a
different pumping. In equation (7.1) for such parameters, one recognizes the Schrödinger equation

$$
\begin{equation*}
i \frac{\partial}{\partial t}|\widetilde{\psi}\rangle=\left[\alpha \widehat{a} \widehat{a}^{\dagger}+V \widehat{a}^{2}+V^{*} \widehat{a}^{\dagger 2}\right]|\widetilde{\psi}\rangle \tag{7.21}
\end{equation*}
$$

for the linear combinations $|\widetilde{\psi}\rangle=\sum \widetilde{\psi}_{n}|2 n\rangle$ of even numbered states of the harmonic oscillator, written with the help of the creation $\widehat{a} \dagger$ and annihilation $\widehat{a}$ operators squared. The odd functions also satisfy the same equation. By analogy to the previous two cases, these squares $\widehat{a}^{2}$ and $\widehat{a}^{\dagger 2}$, together with the unperturbed Hamiltonian $\widehat{a} \dagger \widehat{a}$, also form a closed subalgebra in the entire Hilbert space of the system, satisfying the commutation relations $\left[\widehat{a}^{\dagger} \widehat{a}, \widehat{a}^{2}\right]=$ $-2 \widehat{a}^{2},\left[\widehat{a}^{\dagger} \widehat{a}, \widehat{a}^{\dagger 2}\right]=2 \widehat{a}^{\dagger 2}$, and $\left[\widehat{a}^{2}, \widehat{a}^{\dagger 2}\right]=2 \widehat{a} \widehat{a}^{\dagger}$. Therefore one can try to find an explicit solution for this problem as well.

In the coordinate and momentum representation (7.21) reads

$$
\begin{equation*}
i \frac{\partial}{\partial t}|\widetilde{\psi}\rangle=\left[\alpha \frac{\widehat{Q}^{2}+\widehat{P}^{2}}{2}+\operatorname{Re} V\left(\widehat{Q}^{2}-\widehat{P}^{2}\right)-\operatorname{Im} V(\widehat{P} \widehat{Q}+\widehat{Q} \widehat{P})\right]|\widetilde{\psi}\rangle \tag{7.22}
\end{equation*}
$$

and we note that the operators $\widehat{\mu}_{+}=\left(\widehat{Q}^{2}+\widehat{P}^{2}\right) / 2, \widehat{\mu}_{-}=\left(\widehat{Q}^{2}-\widehat{P}^{2}\right) / 2$, and $\widehat{\mu}_{s}=(\widehat{P} \widehat{Q}+\widehat{Q} \widehat{P}) / 2$ satisfy the commutation relation resembling that for the spin operators $\sigma_{x}, \sigma_{y}$, and $\sigma_{z}$,

$$
\begin{equation*}
\left[\widehat{\mu}_{s}, \widehat{\mu}_{ \pm}\right]=-2 i \widehat{\mu}_{\mp} ; \quad\left[\widehat{\mu}_{+}, \widehat{\mu}_{-}\right]=-2 i \widehat{\mu}_{s} \tag{7.23}
\end{equation*}
$$

which however have an important difference in sign, resulting in a noncompact $S U(1,1)$ Lie group which apart from rotation in the oscillator phase space $\exp \left(-i \phi \widehat{\mu}_{+}\right)$, also contains a squeezing operation $\exp \left(-i Z \widehat{\mu}_{s}\right)$.

Let us assume $V$ real (which can always be done by proper choice of the representation) and perform two sequential transformations $\left|\psi^{\prime}\right\rangle=$ $\exp \left(-i Z \widehat{\mu}_{s}\right) \exp \left(-i \phi \widehat{\mu}_{+}\right)|\widetilde{\psi}\rangle$ with time-dependent $Z$ and $\phi$. The first transformation mixes up the operators $\widehat{\mu}_{-}$and $\widehat{\mu}_{s}$

$$
\begin{align*}
\exp \left(i \phi \widehat{\mu}_{+}\right) \widehat{\mu}_{-} \exp \left(-i \phi \widehat{\mu}_{+}\right) & =\widehat{\mu}_{-} \cos 2 \phi-\widehat{\mu}_{s} \sin 2 \phi \\
\exp \left(i \phi \widehat{\mu}_{+}\right) \widehat{\mu}_{s} \exp \left(-i \phi \widehat{\mu}_{+}\right) & =\widehat{\mu}_{s} \cos 2 \phi+\widehat{\mu}_{-} \sin 2 \phi \tag{7.24}
\end{align*}
$$

in the same way as the transformation $\exp \left(-i \phi \widehat{\sigma}_{x}\right)$ mixes up the $\sigma_{y}$, and $\sigma_{z}$ operators. The second transformation

$$
\begin{align*}
& \exp \left(i Z \widehat{\mu}_{s}\right) \widehat{\mu}_{+} \exp \left(-i Z \widehat{\mu}_{s}\right)=\widehat{\mu}_{+} \cosh 2 Z+\widehat{\mu}_{-} \sinh 2 Z \\
& \exp \left(i Z \widehat{\mu}_{s}\right) \widehat{\mu}_{-} \exp \left(-i Z \widehat{\mu}_{s}\right)=\widehat{\mu}_{-} \cosh 2 Z+\widehat{\mu}_{+} \sinh 2 Z \tag{7.25}
\end{align*}
$$

corresponds to squeezing of the phase space in the coordinate direction and stretching it in the momentum direction:

$$
\begin{align*}
& \exp \left(i Z \widehat{\mu}_{s}\right) \widehat{Q} \exp \left(-i Z \widehat{\mu}_{s}\right)=\widehat{Q} e^{-Z} \\
& \exp \left(i Z \widehat{\mu}_{s}\right) \widehat{P} \exp \left(-i Z \widehat{\mu}_{s}\right)=\widehat{P} e^{Z} \tag{7.26}
\end{align*}
$$

In this representation we find

$$
\begin{align*}
i \frac{\partial}{\partial t}\left|\psi^{\prime}\right\rangle= & {[(\alpha+\dot{\phi}) \cosh 2 Z+2 V \cos \phi \sinh 2 Z] \widehat{\mu}_{+}\left|\psi^{\prime}\right\rangle } \\
& +\left[V \cos \phi \cosh 2 Z+\frac{1}{2}(\alpha+\dot{\phi}) \sinh 2 Z\right] 2 \widehat{\mu}_{-}\left|\psi^{\prime}\right\rangle \\
& -(V \sin \phi-\dot{Z}) 2 \widehat{\mu}_{s}\left|\psi^{\prime}\right\rangle \tag{7.27}
\end{align*}
$$

We now require the coefficients in front of $\widehat{\mu}_{-}$and $\widehat{\mu}_{s}$ to vanish, and determine $Z(t)$ and $\phi(t)$ by solving the set of two coupled nonlinear equations

$$
\begin{align*}
& \dot{\phi}=-\alpha-2 V \cos \phi \operatorname{coth} 2 Z \\
& \dot{Z}=V \sin \phi \tag{7.28}
\end{align*}
$$

which depend on the particular time dependence of $V(t)$ and $\alpha(t)$. The remaining terms in (7.27) result in a transformation $\exp \left(-i \theta \widehat{\mu}_{+}\right)$where the phase reads $\theta=\int[(\alpha+\dot{\phi}) \cosh 2 Z+2 V \cos \phi \sinh 2 Z] d t$, and hence each number state $|n\rangle$ acquires an individual phase shift $n \theta$. The resulting distribution is given by multiplication of the initial condition by the resulting evolution matrix

$$
\begin{equation*}
\left.\left|\psi_{n}\right|^{2}=\left|\sum_{n^{\prime}}\langle n| \exp \left(-i \theta \widehat{\mu}_{+}\right) \exp \left(-i Z \widehat{\mu}_{s}\right) \exp \left(-i \phi \widehat{\mu}_{+}\right)\right| n^{\prime}\right\rangle\left.\psi_{n^{\prime}}\right|^{2} \tag{7.29}
\end{equation*}
$$

and for a single initially populated level $n^{\prime}$, the value of $Z(t)$ is the only parameter governing the distribution at time $t$.

One again finds the population distribution (7.29), considering the overlap integrals of the properly normalized squeezed initial coordinate wavefunction $\exp \left(-i Z \widehat{\mu}_{s}\right) H_{n^{\prime}}(x) e^{-x^{2} / 2}=H_{n^{\prime}}\left(x e^{-Z}\right) e^{-x^{2} e^{-2 Z} / 2}$ and the final wavefunction

$$
\begin{equation*}
\left|\psi_{n}\right|^{2}=\frac{e^{-Z}}{2^{n+n^{\prime}} n!n^{\prime}!}\left|\int H_{n}(x) e^{-x^{2} / 2} H_{n^{\prime}}\left(x e^{-Z}\right) e^{-x^{2} e^{-2 Z} / 2} d x\right|^{2} \tag{7.30}
\end{equation*}
$$

One can say that the symmetry corresponding to the algebraic operator relations (7.23) manifests itself in the coordinate representation - the wavefunction of the oscillator remains the same throughout the evolution apart from the fact that the coordinate itself experiences a scaling by a time-dependent complex number. In particular, this means that Gaussian wavefunctions remain gaussian although the complex width parameter of the Gaussian distribution changes with the course of time. Such wavefunctions are known as squeezed states of the harmonic oscillator.

### 7.1.4 The Harmonic Oscillator in the Simultaneous Presence of Dipole and Raman Pumping

We note that for the harmonic oscillator, the more general problem of simultaneous excitation via both dipole and Raman interaction has an explicit analytical solution. For this case, the closed subalgebra also includes the operators $\widehat{P}^{2}$ and $\widehat{Q}^{2}$, or the operators $\widehat{a}^{\dagger 2}$ and $\widehat{a}^{2}$ in the representation of second quantization. Instead of (7.2), this problem gives a five-term recurrent relation, which being considered with the help of generating functions results in awkward intermediate expressions. That is why we do not present this approach here. However, because of the important role played by harmonic oscillator models in quantum physics, we give here an alternative description.

Consider the Hamiltonian

$$
\begin{equation*}
\widehat{H}=a(t) \widehat{P}^{2}+b(t) \widehat{Q}^{2}+c(t) \widehat{P} \widehat{Q}+f(t) \widehat{P}+h(t) \widehat{Q} \tag{7.31}
\end{equation*}
$$

and perform several sequential transformations in order to find a representation where this expression has the most simple form. We first consider the shift transformation $\widehat{U}_{1}(t)=\exp [i(A(t) \widehat{P}+B(t) \widehat{Q})]$ and find the timedependent shifts of the coordinate $A(t)$ and of the momentum $-B(t)$ such that the transformed Hamiltonian

$$
\begin{equation*}
\widehat{H}_{1}=\widehat{U}_{1}(t) \widehat{H} \widehat{U}_{1}^{-1}(t) \tag{7.32}
\end{equation*}
$$

does not contain terms linear in $\widehat{P}$ and $\widehat{Q}$. To this end we note that

$$
\begin{align*}
& \widehat{U}_{1}(t) \widehat{P} \widehat{U}_{1}^{-1}(t)=\widehat{P}-B(t) \\
& \widehat{U}_{1}(t) \widehat{Q} \widehat{U}_{1}^{-1}(t)=\widehat{Q}+A(t) \tag{7.33}
\end{align*}
$$

as follows from the commutation relation and (7.18), and substitute (7.33) into (7.31), (7.32). This yields

$$
\begin{align*}
\widehat{H}_{1}= & a(t) \widehat{P}^{2}+b(t) \widehat{Q}^{2}+c(t) \widehat{P} \widehat{Q}+[c(t) A(t)-2 a(t) B(t)+f(t)] \widehat{P} \\
& +[2 b(t) A(t)-c(t) B(t)+h(t)] \widehat{Q}+b(t) A(t)^{2} \\
& -c(t) B(t) A(t)+a(t) B(t)^{2}-f(t) B(t)+h(t) A(t), \tag{7.34}
\end{align*}
$$

and for

$$
\begin{equation*}
A(t)=\frac{2 a(t) h(t)-c(t) f(t)}{c^{2}(t)-4 a(t) b(t)} ; B(t)=\frac{2 b(t) f(t)-c(t) h(t)}{c^{2}(t)-4 a(t) b(t)} \tag{7.35}
\end{equation*}
$$

the terms linear in $\widehat{P}$ and $\widehat{Q}$ disappear, whereas the free terms do not contain an operator part and therefore can be ignored.

We have therefore reduced the problem to the former one, and hence the total evolution can be given in the form

$$
\begin{align*}
\widehat{U}(t)= & \exp \left[-\frac{\theta(t)}{2 i}\left(\widehat{P}^{2}+\widehat{Q}^{2}\right)\right] \exp [-i Z(t) \widehat{P} \widehat{Q}] \\
& \exp \left[-i \frac{\phi(t)}{2}\left(\widehat{P}^{2}+\widehat{Q}^{2}\right)\right] \exp [-i A(t) \widehat{P}-i B(t) \widehat{Q}] \tag{7.36}
\end{align*}
$$

as follows from (7.29) whereas five time-dependent coefficients $A(t), B(t)$, $\phi(t), Z(t), \theta(t)$ are related to the arbitrary functions $a(t), b(t), c(t), f(t)$, and $h(t)$ by (7.28), (7.35). For an initial state, localized on just a single number state $|n\rangle$, the population distribution depends only on three parameters, since in this case the operators with $\widehat{P}^{2}+\widehat{Q}^{2}$ in the exponents give only an irrelevant phase dependence, whereas the important part reads

$$
\begin{equation*}
\widehat{U}(t)=\exp [-i Z(t) \widehat{P} \widehat{Q}] \exp \left[-i A^{\prime}(t) \widehat{P}-i B^{\prime}(t) \widehat{Q}\right] \tag{7.37}
\end{equation*}
$$

with

$$
\begin{align*}
A^{\prime}(t) & =A \cos \phi(t)+B \sin \phi(t) \\
B^{\prime}(t) & =B \cos \phi(t)-A \sin \phi(t) \tag{7.38}
\end{align*}
$$

We note that the last exponent in (7.36) corresponds to the case of dipole excitation (7.20) whereas the first exponent is related to the Raman pumping (7.30).

### 7.2 General Case of an Exactly Soluble Relay

We now consider the general case of (7.1) with time independent coefficients, and try to find the conditions where the solution can be expressed explicitly in terms of known functions. The matrix form of the Schrödinger equation

$$
\left\|\left\|\begin{array}{c}
\dot{\psi}_{1}  \tag{7.39}\\
\dot{\psi}_{2} \\
\dot{\psi}_{3} \\
\cdot \\
\dot{\psi}_{N-2} \\
\dot{\psi}_{N-1} \\
\dot{\psi}_{N}
\end{array}\right\|=\right\| \begin{array}{ccccccc}
E_{1} & V_{1}^{*} & 0 & . & 0 & 0 & 0 \\
V_{1} & E_{2} & V_{2}^{*} . & 0 & 0 & 0 \\
0 & V_{2} & E_{3} & . & 0 & 0 & 0 \\
. & \cdot & \cdot & \cdot & . & . & . \\
0 & 0 & 0 & . & E_{N-2} & V_{N-2}^{*} & 0 \\
0 & 0 & 0 & . & V_{N-2} & E_{N-1} & V_{N-1}^{*} \\
0 & 0 & 0 & . & 0 & V_{N-1} & E_{N}
\end{array}\left\|\begin{array}{c}
\psi_{1} \\
\psi_{2} \\
\psi_{3} \\
. \\
\psi_{N-2} \\
\psi_{N-1} \\
\psi_{N}
\end{array}\right\|,
$$

allows one to write, for the Fourier transforms

$$
\left\|\begin{array}{c}
\psi_{1}  \tag{7.40}\\
\psi_{2} \\
\cdot \\
\psi_{N-1} \\
\psi_{N}
\end{array}\right\|=i\left\|\begin{array}{ccccc}
E_{1}-\varepsilon & V_{1}^{*} & \cdot & 0 & 0 \\
V_{1} & E_{2}-\varepsilon & . & 0 & 0 \\
\cdot & \cdot & \cdot & \cdot & \cdot \\
0 & 0 & \cdot & E_{N-1}-\varepsilon & V_{N-1}^{*} \\
0 & 0 & \cdot & V_{N-1} & E_{N}-\varepsilon
\end{array}\right\|^{-1}\left\|\begin{array}{c}
\psi_{1}^{\prime} \\
\psi_{2}^{\prime} \\
\cdot \\
\psi_{N-1}^{\prime} \\
\psi_{N}^{\prime}
\end{array}\right\|
$$

where primes denote initial conditions at $t=0$. In accordance with the standard rules of linear algebra, the inverse of the three-diagonal matrix $\widehat{H}-\varepsilon$ can be found with the help of the diagonal minors, $M_{k}^{+}$corresponding to the $k \times k$ upper-left corner and $M_{n}^{-}$corresponding to the $(N-n) \times(N-n)$ lower-right corner of the matrix $\widehat{H}-\varepsilon$, and for the matrix elements with $k<n$ one finds

$$
\begin{equation*}
(\widehat{H}-\varepsilon)_{k n}^{-1}=\frac{M_{k}^{+} M_{n}^{-}}{\operatorname{Det}|\widehat{H}-\varepsilon|} \prod_{l=k}^{n-1} V_{l}^{*} \tag{7.41}
\end{equation*}
$$

that for the matrix element with $k=4$ and $n=N-3$ in particular means

$$
(\widehat{H}-\varepsilon)_{4, N-3}^{-1}=\frac{\left|\begin{array}{ccc}
E_{1} & V_{1}^{*} & 0  \tag{7.42}\\
V_{1} & E_{2} & V_{2}^{*} \\
0 & V_{2} & E_{3}
\end{array}\right|\left|\begin{array}{ccc}
E_{N-2} & V_{N-2}^{*} & 0 \\
V_{N-2} & E_{N-1} & V_{N-1}^{*} \\
0 & V_{N-1} & E_{N}
\end{array}\right|}{\operatorname{Det}|\widehat{H}-\varepsilon|} V_{4}^{*} \ldots V_{N-4}^{*}
$$

The minors $M_{n}^{-}$and $M_{k}^{+}$satisfy the three-term recurrence relations

$$
\begin{align*}
& M_{n}^{-}=\left(E_{n+1}-\varepsilon\right) M_{n+1}^{-}-\left|V_{n+1}\right|^{2} M_{n+2}^{-} \\
& M_{k}^{+}=\left(E_{k-1}-\varepsilon\right) M_{k-1}^{+}-\left|V_{k-1}\right|^{2} M_{k-2}^{+} \tag{7.43}
\end{align*}
$$

which coincide with the Fourier transformed initial Schrödinger equation. Indeed, after the replacement $\psi_{n} \rightarrow \phi_{n} \prod_{k=0}^{n-1} V_{k}^{*}$, equation (7.39) yields

$$
\begin{equation*}
\left(E_{n}-\varepsilon\right) \phi_{n}=\phi_{n-1}+\left|V_{n}\right|^{2} \phi_{n+1} \tag{7.44}
\end{equation*}
$$

where we have ignored for the moment the initial conditions in the form of $\delta$-functions, discussed earlier in the context of (3.9). We also note that the relation

$$
\begin{equation*}
M_{n+1}^{+} M_{n}^{-}-\left|V_{n}\right|^{2} M_{n}^{+} M_{n+1}^{-}=\operatorname{Det}|\widehat{H}-\varepsilon| \tag{7.45}
\end{equation*}
$$

is valid for any $n$.

### 7.2.1 Conditions for the Existence of a Polynomial Solution

Let us consider the conditions where the recurrence relations (7.43), (7.44) have an explicit algebraic solution, representable by an analytic expression. We concentrate only on the particular case where $E_{n+1}$ is a polynomial in $n$ and $\left|V_{n+1}\right|^{2}=F_{n} / D_{n}$ is a ratio of two polynomials. The minors in this case will also be ratios of polynomials, as well as the combinations $u_{n}^{-}=$ $M_{n}^{-} / M_{n+1}^{-}$and $u_{k}^{+}=M_{k}^{+} / M_{k-1}^{+}$satisfying the equations

$$
\begin{align*}
& u_{n}^{-}=\left(E_{n+1}-\varepsilon\right)-\frac{\left|V_{n+1}\right|^{2}}{u_{n+1}^{-}} \\
& u_{k}^{+}=\left(E_{k-1}-\varepsilon\right)-\frac{\left|V_{k-1}\right|^{2}}{u_{k-1}^{+}} \tag{7.46}
\end{align*}
$$

These two equations are equivalent. Therefore, we take just the second one and substitute $u_{k}^{+}=P_{k} / Q_{k}$ in the form of a ratio of a polynomial $P_{k}$ of order $p$ and a polynomial $Q_{k}$ of order $q$. This yields

$$
\begin{equation*}
D_{k} P_{k}^{+} P_{k+1}^{+}=D_{k}\left(E_{k}-\varepsilon\right) P_{k}^{+} Q_{k+1}^{+}-F_{k} Q_{k}^{+} Q_{k+1}^{+} \tag{7.47}
\end{equation*}
$$

Now one can ask if a solution of (7.47) exists for $p$ and $q$ independent of $k$. In this case, the minors can be found in the form of an explicit algebraic expression. Otherwise no general algebraic expression exists for the minors simpler than their definitions. The $p+q+2$ coefficients of the polynomials $P_{k}$ and $Q_{k}$ play the role of independent variables in (7.47), whereas the coincidence of the coefficients corresponding to the same powers of $k$ on the left and on the right hand sides results in a set of at least $2 p+d+1$ equations, where $d$ is the order of the polynomial $D_{k}$, as suggested by the polynomial order of left hand side of (7.47). A solution exists for $p+q+2 \geq 2 p+d+1$, that is for $q-p \geq d-1$. Moreover, the last term on the right hand side suggests that the number of equations is not less than $2 q+f+1$, where $f$ is the order of $F_{k}$, and therefore $d-1 \leq q-p \leq 1-f$. Therefore $q-p$ can only assume the values $0, \pm 1$. One more condition $1 \geq d+\epsilon$ emerges from the polynomial order of the first term on the right-hand side, where $\epsilon$ is the polynomial order of $E_{k}$. The latter yields two possibilities: $d=1 ; \epsilon=0$, or $d=0 ; \epsilon=1$, that imply $f \leq 1$ for the first and $f \leq 2$ for the second cases.

The case $d=0 ; \epsilon=1 ; f=0$ coincides with (7.3), and the case $d=0 ; \epsilon=$ $1 ; f=1$ corresponds to the harmonic oscillator (7.12). Hence we are left with only two more possibilities

$$
\begin{array}{ll}
d=1 ; & \epsilon=0 ; \quad f=1 ; \quad(p=q) \\
d=0 ; \quad \epsilon=1 ; \quad f=2 ; \quad(p=q+1) \tag{7.48}
\end{array}
$$

where $d=1, \epsilon=0, f=0$ is a particular case of the first possibility. The second possibility includes Raman pumping (7.21) as a particular case. Equations (7.43) for minors now read

$$
\begin{align*}
& (a n+b)\left(M_{n-1}^{-}+\varepsilon M_{n}^{-}\right)+M_{n+1}^{-}=0 \\
& (a k+b)\left(M_{k+1}^{+}+\varepsilon M_{k}^{+}\right)+M_{k-1}^{+}=0 \tag{7.49}
\end{align*}
$$

for the decreasing coupling $\left|V_{n}\right|^{-2}=a n+b$; and strict resonance $E_{n}=0$, and

$$
\begin{align*}
M_{n-1}^{-} & =(\alpha n-\varepsilon) M_{n}^{-}-a(n-b)(n-c) M_{n+1}^{-} \\
M_{k+1}^{+} & =(\alpha k-\varepsilon) M_{k}^{+}-a(k-b)(k-c) M_{k-1}^{+} \tag{7.50}
\end{align*}
$$

for the increasing coupling $\left|V_{n}\right|^{2}=a(n-b)(n-c)$ and an arbitrary linear detuning $E_{n}=\alpha n$. These couplings are shown in Fig. 7.2. The corresponding equations (7.44) read


Fig. 7.2. Two types of dependences of the coupling strength on the level number corresponding to exactly soluble problems (Sect. 7.2.2, left and Sect. 7.2.3, right). Arrows show the domains of the imaginary matrix elements.

$$
\begin{align*}
& (a n+b)\left(\varepsilon \phi_{n}+\phi_{n-1}\right)+\phi_{n+1}=0 \\
& (\alpha n-\varepsilon) \phi_{n}-\phi_{n-1}-a(n-b)(n-c) \phi_{n+1}=0 \tag{7.51}
\end{align*}
$$

For the case of increasing coupling, it is convenient to make the replacement $\phi_{n} \rightarrow \phi_{n} \Gamma(1-c) / \Gamma(n-c)$. Equations (7.51), with allowance of the initial conditions $\psi_{n}(0)$ at $t=0$, ignored earlier in (7.44), take the form

$$
\begin{align*}
& (a n+b)\left(\varepsilon \phi_{n}+\phi_{n-1}\right)+\phi_{n+1}=i(a n+b) \psi_{n}(0) \prod_{k=0}^{n-1}\left(V_{k}^{*}\right)^{-1}, \\
& (\varepsilon-\alpha n) \phi_{n}+(n-1-c) \phi_{n-1}+a(n-b) \phi_{n+1}=i \psi_{n}(0) \prod_{k=0}^{n-1} \frac{(k-c)}{V_{k}^{*}} . \tag{7.52}
\end{align*}
$$

Equations (7.52) are linear in $n$, and can be solved as earlier (7.5), by introducing the generating functions $\mathcal{M}(\varphi)=\sum_{n} \phi_{n} e^{i \varphi n}$. The generating function method suggests the replacement $n \rightarrow-i \partial / \partial \varphi$, and we arrive at the first-order differential equations

$$
\begin{align*}
& i a \frac{\partial\left(e^{i \varphi}+\varepsilon\right) \mathcal{M}(\varphi)}{\theta \varphi}=\left(b e^{i \varphi}+b \varepsilon+e^{-i \varphi}\right) \mathcal{M}(\varphi)+q(\varphi)  \tag{7.53}\\
& i \frac{\partial\left(\alpha-a e^{-i \varphi}-e^{i \varphi}\right) \mathcal{M}(\varphi)}{\theta \varphi}=\left[(1+c) e^{i \varphi}-\varepsilon+a b e^{-i \varphi}\right] \mathcal{M}(\varphi)+q(\varphi), \tag{7.54}
\end{align*}
$$

with the inhomogeneous terms $q(\varphi)$ originating from the initial conditions.
Note that .(7.53)-(7.54) are first-order linear differential equations, and that they always have explicit solutions in the form of integrals. The reason why the original problem reduces to an exactly soluble one originates in the condition (7.47), which implies a sort of translational symmetry, which follows from the requirement of the independence of the polynomial orders of $P_{k}$ and $Q_{k}$ entering (7.46) for $u_{k}^{+}=P_{k} / Q_{k}$ on the index $k$. Moreover, it points
to the fact that the Schrödinger equation, with the detunings and interactions restricted by the condition (7.48), describe evolution in a subspace of the entire Hilbert space of the multilevel system, whereas the corresponding evolution operators belong to a certain subgroup. For certain particular cases of increasing coupling, this is the subgroup $S U(2)$ of rotations generated by the angular momentum operators $L_{ \pm}, L_{z}$, whereas the Hamiltonian can be written as a sum $a L_{+}+a^{*} L_{-}+\alpha L_{z}$ of these operators.

The homogeneous parts of (7.53)-(7.54) have the form

$$
\begin{equation*}
i \frac{\partial f(\varphi) \mathcal{M}(\varphi)}{\theta \varphi}=g(\varphi) \mathcal{M}(\varphi) \tag{7.55}
\end{equation*}
$$

where $f(x)$ and $g(x)$ are arbitrary functions. Therefore their solution has the form

$$
\begin{equation*}
\mathcal{M}(\varphi)=\frac{1}{f(\varphi)} \exp \left\{-i \int^{\varphi} \frac{g(x)}{f(x)} d x\right\} \tag{7.56}
\end{equation*}
$$

which results in a general solution of the homogenous equation

$$
\begin{equation*}
\phi_{n}=C \int_{0}^{2 \pi} \frac{e^{-i \varphi n}}{f(\varphi)} \exp \left\{-i \int^{\varphi} \frac{g(x)}{f(x)} d x\right\} d \varphi \tag{7.57}
\end{equation*}
$$

The complete equations

$$
\begin{equation*}
i \frac{\partial f(\varphi) \mathcal{M}(\varphi)}{\theta \varphi}=g(\varphi) \mathcal{M}(\varphi)+q(\varphi) \tag{7.58}
\end{equation*}
$$

have the solution

$$
\begin{equation*}
\mathcal{M}(\varphi)=\frac{1}{i f(\varphi)} \int^{\varphi} q(y) \exp \left\{i \int_{\varphi}^{y} \frac{g(x)}{f(x)} d x\right\} d y \tag{7.59}
\end{equation*}
$$

and hence

$$
\begin{equation*}
\phi_{n}=\int_{0}^{2 \pi} \frac{e^{-i \varphi n}}{i f(\varphi)}\left[\int^{\varphi} q(y) \exp \left\{i \int_{\varphi}^{y} \frac{g(x)}{f(x)} d x\right\} d y\right] d \varphi \tag{7.60}
\end{equation*}
$$

is the solution of the original problem. Straightforward calculations of the integrals (7.60) are rather awkward. We therefore separately consider the case of increasing coupling, for which another approach is more convenient, and then we concentrate on the case of decreasing coupling.

### 7.2.2 The Increasing Coupling $\left|V_{n}\right|=\sqrt{a(n-b)(n-c)}$.

Equation (7.54), after inverse Fourier transformation, reads

$$
\begin{equation*}
i\left(\alpha-a e^{-i \varphi}-e^{i \varphi}\right) \frac{\partial \mathcal{M}(\varphi)}{\partial \varphi}+i \frac{\partial \mathcal{M}(\varphi)}{\partial t}=\left[c e^{i \varphi}+a(b-1) e^{-i \varphi}\right] \mathcal{M}(\varphi) \tag{7.61}
\end{equation*}
$$

where the factor $a$ can be set to one, by introducing $\sqrt{a}$ as the frequency unit. As any of the first order linear differential equations, it has a general solution

$$
\begin{equation*}
\mathcal{M}(\varphi, t)=e^{i \varphi(2+b)}\left[\left(e^{i \varphi}-x_{-}\right)\left(e^{i \varphi}-x_{+}\right)\right]^{A_{2}}\left(\frac{e^{i \varphi}-x_{+}}{e^{i \varphi}-x_{-}}\right)^{A_{1}} \mathcal{F}\left(\frac{e^{i \varphi}-x_{-}}{e^{i \varphi}-x_{+}} e^{t \sqrt{4 a+\alpha^{2}}}\right), \tag{7.62}
\end{equation*}
$$

found by the method of characteristics, where

$$
\begin{equation*}
x_{ \pm}=\frac{i \alpha}{2} \pm \frac{i \sqrt{4+\alpha^{2}}}{2} ; A_{1}=i \frac{1+b+c}{2 \sqrt{4-\alpha^{2}}} \alpha ; A_{2}=\frac{c-3-b}{2} \tag{7.63}
\end{equation*}
$$

and $\mathcal{F}(x)$ is an arbitrary function.
The choice of $\mathcal{F}(x)$ is suggested by the initial condition $\mathcal{M}(\varphi, 0)=$ $\sum \phi_{n}(t=0) e^{i \varphi n}$

$$
\begin{equation*}
\mathcal{F}\left(\frac{e^{i \varphi}-x_{-}}{e^{i \varphi}-x_{+}}\right)=e^{-i \varphi(2+b)}\left(\frac{e^{i \varphi}-x_{-}}{e^{i \varphi}-x_{+}}\right)^{A_{1}} \frac{\mathcal{M}(\varphi, 0)}{\left[\left(e^{i \varphi}-x_{-}\right)\left(e^{i \varphi}-x_{+}\right)\right]^{A_{2}}} \tag{7.64}
\end{equation*}
$$

which after the substitution $e^{i \varphi}=\left(y x_{-}-x_{+}\right) /(y-1)$ yields

$$
\begin{equation*}
\mathcal{F}(y)=\frac{y^{-A_{2}-A_{1}}(y-1)^{c-1}}{\left(y x_{-}-x_{+}\right)^{b+2}\left(x_{+}-x_{-}\right)^{2 A_{2}}} \mathcal{M}\left(-i \ln \left(\frac{y x_{-}-x_{+}}{y-1}\right), 0\right) \tag{7.65}
\end{equation*}
$$

Since any initial condition $\phi_{n}(t=0)=\delta_{n, n_{0}}$ after the replacements $b \rightarrow$ $b+n_{0}, c \rightarrow c+n_{0}$ results in $\phi_{n}(t=0)=\delta_{n, 0}$, we can consider only $\mathcal{M}(\varphi, 0)=$ 1 and obtain

$$
\begin{equation*}
\phi_{k}(t)=\int_{0}^{2 \pi} \frac{d \varphi}{2 \pi} \frac{e^{\left(A_{2}-A_{1}\right) T-i \varphi(k-b-2)}}{\left(x_{+}-x_{-}\right)^{2 A_{2}}} \frac{\left(e^{i \varphi}\left(e^{T}-1\right)-x_{+} e^{T}+x_{-}\right)^{c-1}}{\left(e^{i \varphi}\left(e^{T} x_{-}-x_{+}\right)+1-e^{T}\right)^{b+2}} \tag{7.66}
\end{equation*}
$$

where $T=t \sqrt{4 a+\alpha^{2}}$. The other form of the same expression

$$
\begin{align*}
\phi_{k}(t)= & \int_{0}^{2 \pi} d \varphi \frac{e^{-i \varphi k}}{2 \pi} \frac{e^{\left(A_{2}-A_{1}\right) T}}{\left(x_{+}-x_{-}\right)^{2 A_{2}}} \frac{\left(x_{-}-x_{+} e^{T}\right)^{c-1}}{\left(e^{T} x_{-}-x_{+}\right)^{b+2}} \\
& \left(\frac{e^{i \varphi}\left(e^{T}-1\right)}{-x_{+} e^{T}+x_{-}}+1\right)^{c-1}\left(1+\frac{1-e^{T}}{e^{i \varphi}\left(e^{T} x_{-}-x_{+}\right)}\right)^{-b-2} \tag{7.67}
\end{align*}
$$

is more convenient for casting in power series over $\left(e^{T}-1\right)$.
This results in

$$
\begin{align*}
\psi_{k}(t)= & \int_{0}^{2 \pi} d \varphi \frac{e^{-i \varphi k}}{2 \pi} \frac{e^{\left(A_{2}-A_{1}\right) T}}{\left(x_{+}-x_{-}\right)^{2 A_{2}}} \frac{\left(x_{-}-x_{+} e^{T}\right)^{c-1}}{\left(e^{T} x_{-}-x_{+}\right)^{b+2}} \\
& \sum_{m=0}^{\infty} \frac{\Gamma(c) e^{i m \varphi}}{\Gamma(c-m) m!}\left(\frac{e^{T}-1}{x_{-}-x_{+} e^{T}}\right)^{m} \\
& \sum_{n=0}^{\infty} \frac{\Gamma(-1-b) e^{-i n \varphi}}{\Gamma(-1-b-n) n!}\left(\frac{1-e^{T}}{e^{T} x_{-}-x_{+}}\right)^{n} \tag{7.68}
\end{align*}
$$

where we have assumed that $b<0, c>0$. Integration yields $m=k+n$ and we are left with the known sum

$$
\begin{equation*}
\sum_{m=0}^{\infty} \frac{\Gamma(c)}{\Gamma(c-k-n)(-k-n)!} \frac{\Gamma(-1-b)}{\Gamma(-1-b-n) n!} x^{n}=\frac{\Gamma(c)}{k!\Gamma(c-k)}{ }_{2} F_{1}(2+b, 1+k-c ; k+1 ; x), \tag{7.69}
\end{equation*}
$$

given in terms of hypergeometric function ${ }_{2} F_{1}(\alpha, \beta ; \gamma ; z)$. Substitution of this sum yields for the probability amplitudes

$$
\begin{align*}
\psi_{k}(t)= & \frac{i^{k} e^{\left(A_{2}-A_{1}\right) T}}{\left(x_{+}-x_{-}\right)^{2 A_{2}}} \sqrt{\frac{\Gamma(c) \Gamma(k-b)}{\Gamma(c-k) \Gamma(-b)}} \frac{\left(x_{-}-x_{+} e^{T}\right)^{c-1-k}}{\left(e^{T} x_{-}-x_{+}\right)^{b+2}} \frac{\left(e^{T}-1\right)^{k}}{k!} \\
& { }_{2} F_{1}\left(2+b, 1+k-c ; k+1 ; \frac{e^{T}-1}{x_{-}-x_{+} e^{T}} \frac{1-e^{T}}{e^{T} x_{-}-x_{+}}\right) \tag{7.70}
\end{align*}
$$

Note that the transformation formulas for the hypergeometric function can give an expression, symmetric with respect to $b$ and $c$, such that the imposed conditions $c>b$ can be lifted. Otherwise one has to make the replacements $b \rightleftharpoons c, x_{-} \rightleftharpoons x_{+}$. In the opposite case $\alpha<1$, the population distribution remains localized near the initially populated state. Note that in order to get the correct result of the form (7.70) in the domain where $\left|V_{n}\right| \sim \sqrt{(n-b)(n-c)}$ assumes imaginary values, one has to make the replacements $\alpha \rightarrow i \alpha$ and $t \rightarrow i t$. Also note that for the particular case of integer $b$ and $c$, the problem describes the dynamics of an angular momentum $J$, with $2 J+1=|b-c|$, in a constant magnetic field, and therefore corresponds to a $2 L+1$-dimensional representation of the $S U_{2}$ group.

In Fig. 7.3 we show two different regimes of the population dynamics corresponding to small detunings $\alpha<1$ and large detunings $\alpha>1$. In the first regime, the population initially localized at the level $k=0$ moves toward infinite values of $k$ in the direction of increasing couplings $V_{n}$. We note that many different wavepackets are formed during the time evolution, and these packets follow different trajectories.


Fig. 7.3. Two regimes of population dynamics for a system with linear detuning $E_{n}=\alpha n$ and coupling $V=\sqrt{a(n-b)(n-c)}$. (a) In the regime of small detuning $\alpha<\sqrt{a}$, the population distribution moves with the course of time to the levels with large numbers. (b) For the large detunings $\alpha>\sqrt{a}$ the distribution oscillates around the initially populated state. Arrows show the directions of propagation of different wavepackets comprising the population distribution.

### 7.2.3 Decreasing Coupling $\left|V_{n}\right|=1 / \sqrt{a n+b}$

Consider the case of decreasing coupling. Calculating the integrals entering the exponents of (7.57), (7.59) one finds

$$
\begin{equation*}
\int^{\varphi} \frac{b e^{i x}+b \varepsilon+e^{-i x}}{e^{i x}+\varepsilon} d x=b \varphi+\frac{i e^{-i \varphi}}{\varepsilon}-\frac{i}{\varepsilon^{2}} \ln \left(1+\varepsilon e^{-i \varphi}\right) \tag{7.71}
\end{equation*}
$$

and hence (7.57) yields a solution

$$
\begin{align*}
\phi_{k} & =A \int_{0}^{2 \pi} \frac{e^{-i \varphi(k+b)}}{\left(1+\varepsilon e^{-i \varphi}\right)^{1+1 / \varepsilon^{2}}} \exp \left\{\frac{e^{-i \varphi}}{\varepsilon}\right\} d e^{i \varphi} \\
& =\varepsilon^{k} A^{\prime} \int_{C} \frac{y^{(k+b)}}{(1+y)^{1+1 / \varepsilon^{2}}} \exp \left\{\frac{y}{\varepsilon^{2}}\right\} d e^{i \varphi} \tag{7.72}
\end{align*}
$$

for the homogeneous equation (7.52). Here $A$ and $A^{\prime}$ are constants, and the replacement $e^{i \varphi} \rightarrow y / \varepsilon$ has been performed in the last line. We have also set $a=1$ by employing the scaled time $t \rightarrow t \sqrt{a}$ and frequency $\varepsilon \rightarrow \varepsilon / \sqrt{a}$.

Note that the integration contour $C$ must not necessarily be a circle of radius $\varepsilon$ around the point $y=0$, as suggested by the replacement. On the contrary, any contour, such that the integrand takes identical values at his ends, can be chosen in the spirit of the Laplace contour integral method (5.137). In this case, one can show by integration by parts that (7.72) satisfies the homogeneous recurrence relation (7.52). For $k+b>0$, among the possibilities
is the contour of the type shown in Fig. 3.23(b) which goes around the points $y=0$ and $y=-1$ and yields the confluent hypergeometric function

$$
\begin{align*}
& \chi_{k}= b \frac{\varepsilon^{k-1} e^{-1 / \varepsilon^{2}} \Gamma\left(\frac{1}{\varepsilon^{2}}-b\right)}{\Gamma\left(1-b-k+\frac{1}{\varepsilon^{2}}\right)}{ }_{1} F_{1}\left(-\frac{1}{\varepsilon^{2}}, 1+b-\frac{1}{\varepsilon^{2}} ; \frac{1}{\varepsilon^{2}}\right) \\
&{ }_{1} F_{1}\left(-k-b, 1-k-b+\frac{1}{\varepsilon^{2}} ; \frac{1}{\varepsilon^{2}}\right), \tag{7.73}
\end{align*}
$$

where a convenient normalization has been done.
For $k+b<0$ one finds the other solution by noting that the combination $\tilde{\chi}_{k}(b, \varepsilon)=(-i)^{k} \chi_{-k}(1-b, i \varepsilon) \Gamma(k+b) / \Gamma(b)$ also satisfies the homogeneous recurrent relation (7.52). The substitution yields

$$
\begin{align*}
\chi_{k}^{\prime}= & \frac{(1-b)}{(-\varepsilon)^{k-1}} e^{-1 / \varepsilon^{2}} \frac{\Gamma\left(b-1-\frac{1}{\varepsilon^{2}}\right)}{\Gamma\left(b+k-\frac{1}{\varepsilon^{2}}\right)}{ }_{1} F_{1}\left(\frac{1}{\varepsilon^{2}}, 2-b+\frac{1}{\varepsilon^{2}} ; \frac{1}{\varepsilon^{2}}\right) \\
& \frac{\Gamma(k+b)}{\Gamma(b)} \quad{ }_{1} F_{1}\left(1-\frac{1}{\varepsilon^{2}}, k+b-\frac{1}{\varepsilon^{2}} ; \frac{1}{\varepsilon^{2}}\right), \tag{7.74}
\end{align*}
$$

where the relation ${ }_{1} F_{1}(\alpha, \gamma ; x)=e^{x}{ }_{1} F_{1}(\gamma-\alpha, \gamma ;-x)$ has been employed. The normalization of $\chi_{k}$ and $\chi_{k}^{\prime}$ has been chosen such that the combination

$$
\begin{equation*}
\phi_{k}=\chi_{k} \Theta(k)+\chi_{k}^{\prime} \Theta(-k) \tag{7.75}
\end{equation*}
$$

satisfies (7.52) with $\psi_{0}(t=0)=1$. This enables us to avoid cumbersome calculations of the integrals suggested by the explicit expression (7.60). Timedependent amplitudes have to be found by numerical calculation of the inverse Fourier transformation of (7.75).

It is worth mentioning that $k$-dependence of the solutions (7.73)-(7.74) coincide with $k$-dependence of the minors (7.49), and hence the determinant Det $|\widehat{H}-\varepsilon|$ in (7.41) behaves as $\chi_{k \rightarrow \infty}^{\prime}$. For $|\varepsilon|<1$ it vanishes at the points $\varepsilon=1 / \sqrt{b+n}$ of the complex $\varepsilon$-plane, where $n$ is an integer. This gives the frequency spectrum of the system. Note that $1 / \sqrt{b+n}$ coincides by accident with the Rabi frequencies of two-level systems composed of neighboring levels. In Fig. 7.4 one sees the interplay of these non commensurable eigenfrequencies resulting in a partial transfer of the population to the states with large numbers. The corresponding wavepackets, however, lose their population in the course of propagation at a constant velocity.

### 7.3 Smooth Variation of the Parameters

As we have seen, the exactly soluble problems correspond to certain subgroups in the Hilbert space of the multilevel systems. This implies that the corresponding Hamiltonians in (7.1) consist of several additive parts that


Fig. 7.4. Two regimes of population dynamics for a resonant system with a decreasing coupling $V=1 / \sqrt{a(n-b)}$. One sees that the interference of different non-comensurate harmonics results in involved population dynamics: a part of the population remains at the levels neighboring the initially populated state, whereas the rest of the population leaves toward states with large numbers $n$. The last part of the population distribution forms wavepackets, gradually losing their population with the course of propagation toward large $n$ with a constant velocity. The arrows show the dynamics of the main wavepackets formed as a result of the interference.
form an algebra of relatively small size, that is commutators of these parts are given by their linear combinations. This is not always the case for systems with couplings $V_{n}$ and energies $E_{n}$ arbitrarily changing with the level number $n$. However, if these quantities change slowly, the commutators are small and can be ignored.

### 7.3.1 WKB approximation

In this extreme, one can also consider the level number as a continuous coordinate, and write the Schrödinger equation (7.1) in the form

$$
\begin{equation*}
i \frac{\partial \psi(n)}{\partial t}=E(n) \psi(n)+V(n) \exp \left(-\frac{\partial}{\partial n}\right) \psi(n)+V^{*}(n+1) \exp \left(\frac{\partial}{\partial n}\right) \psi(n) \tag{7.76}
\end{equation*}
$$

where the finite shift operator (5.12) has been employed. The generating function $\Psi(t, p)=\sum_{n} \psi_{n}(t) e^{i p n}$ of (7.5) may now be interpreted as the probability amplitude in the momentum representation, while the finite shift operator in this representation is multiplicative in nature: $\exp \left(\frac{\partial}{\partial n}\right) \rightarrow e^{i p}$.

Equation (7.76) can be interpreted as the Schrödinger equation of a particle with a non-quadratic relation between energy and momentum, moving in direction $n$. The momentum operator $\widehat{p}=-i \frac{\partial}{\partial n}$ commutes neither with
the analog of the potential energy $E(n)$ nor with the analog of the inverse mass $V(n)$, and the commutators given by the derivatives of these functions, generally nonlinear, do not commute with $\widehat{p}$ either. However, when we consider the dimensionless mechanical action $S(n, t)$ given by the WKB substitution $\psi(n, t)=e^{-i S(n, t)}$ for slowly changing $E(n)$ and $V(n)$, the higherorder derivatives turn out to be small, as compared to the powers of the first derivative. In other words, higher-order commutators of the action with the momentum operator can be ignored which yields

$$
\begin{equation*}
\frac{\partial S(n, t)}{\partial t}=E(n)+V(n) \exp \left(i \frac{\partial S(n, t)}{\partial n}\right)+V^{*}(n) \exp \left(\frac{\partial S(n, t)}{i \partial n}\right) \tag{7.77}
\end{equation*}
$$

where the difference between $V^{*}(n+1)$ and $V^{*}(n)$ has also been ignored within the order of the approximation. In this equation, one recognizes the Hamilton-Jacobi equation for a classical particle with the Hamiltonian

$$
\begin{equation*}
\mathcal{H}(n, p)=E(n)+2 V(n) \cos p \tag{7.78}
\end{equation*}
$$

where $2 V(n) \cos p$ plays the role of the kinetic energy. Let us first assume that the classical total energy of the system is $\varepsilon$. Then one immediately sees two regimes: (a) classically forbidden $|E(n)-\varepsilon|>|2 V(n)|$, where the potential energy exceeds the kinetic energy, and (b) classically allowed, where $|E(n)-\varepsilon|<|2 V(n)|$.

Substitution of $S(n, t)=\varepsilon t+S_{m}(n)$ into (7.78) yields the equation

$$
\begin{equation*}
\frac{\partial S_{m}(n)}{\partial n}=\arccos \frac{\varepsilon-E(n)}{2 V(n)} \tag{7.79}
\end{equation*}
$$

for the reduced action $S_{m}(n)$ which determines the eigenfunction $\psi(n, \varepsilon)=$ $e^{-i S_{m}(n)}$ for the eigenstate of energy $\varepsilon$. This allows us to find $S_{m}(n)$ explicitly, and we obtain the WKB approximate result

$$
\begin{equation*}
\psi(n, t)=\int f(\varepsilon) \exp \left[-i \varepsilon t-i \int^{n} \arccos \frac{\varepsilon-E(x)}{2 V(x)} d x\right] d \varepsilon \tag{7.80}
\end{equation*}
$$

which has a clear physical meaning: the wavefunction is a linear superposition of the eigenfunctions of different energies $\varepsilon$ with arbitrary amplitudes $f(\varepsilon)$. The latter have to be found from the initial conditions. Note that in the particular case $E(n)=$ const $; V(n)=$ const, considered on p. 310, the amplitude (7.80) gives an exact solution of the problem.

### 7.3.2 Position and Width of the Erenfest Wavepacket

Let us analyze some physical consequences of (7.80). Let us assume that initially the wave packet is localized near the point $x=x_{0}$ at zero energy. At
long times, the main contribution to the integral over energies $\varepsilon$ comes from the vicinity of the point $\varepsilon=0$. We therefore can develop the classical action $S(n, t)$

$$
\begin{equation*}
S(n, t)=\varepsilon t+\int^{n} \arccos \frac{\varepsilon-E(x)}{2 V(x)} d x=0 \tag{7.81}
\end{equation*}
$$

in a Taylor series near this point. This yields

$$
\begin{align*}
S(n, t)= & \int^{n} \arccos \frac{-E(x)}{2 V(x)} d x  \tag{7.82}\\
& +\varepsilon\left(t-\int^{n} \frac{d x}{\left[4 V(x)^{2}-E(x)^{2}\right]^{1 / 2}}\right)+\frac{\varepsilon^{2}}{4} \int^{n} \frac{E(x) d x}{\left[4 V(x)^{2}-E(x)^{2}\right]^{3 / 2}} .
\end{align*}
$$

A relatively smooth function $f(\varepsilon)$ can be factored out of the integral (7.80) and replaced by its value at zero, and the integration with the allowance of (7.82) yields

$$
\begin{align*}
\psi(n, t)= & \frac{f(0) \sqrt{2 \pi}}{\sqrt{\left[\frac{n}{[ } \frac{E(x) d x}{\left[4 V(x)^{2}-E(x)^{2}\right]^{3 / 2}}\right.}} \exp \left[-i \int^{n} \arccos \frac{-E(x)}{2 V(x)} d x\right] \\
& \exp \left[-\left(t+\int^{n} \frac{d x}{\left[4 V(x)^{2}-E(x)^{2}\right]^{1 / 2}}\right)^{2} / \int^{n} \frac{E(x) d x}{\left[4 V(x)^{2}-E(x)^{2}\right]^{3 / 2}}\right] . \tag{7.83}
\end{align*}
$$

This means that a wavepacket passes the point $n$ at time

$$
\begin{equation*}
t=\int^{n} \frac{d x}{\left[4 V(x)^{2}-E(x)^{2}\right]^{1 / 2}}, \tag{7.84}
\end{equation*}
$$

at velocity

$$
\begin{equation*}
\frac{\partial n}{\partial t}=\left[4 V(n)^{2}-E(n)^{2}\right]^{1 / 2} \tag{7.85}
\end{equation*}
$$

and the passage takes time

$$
\begin{equation*}
\delta t=\int^{n} \frac{E(x) d x}{\left[4 V(x)^{2}-E(x)^{2}\right]^{3 / 2}} \tag{7.86}
\end{equation*}
$$

Note that the width $\delta t$ is given by the integral diverging at the classical turning points where $4 V(x)^{2}=E(x)^{2}$, and the main contribution comes from the vicinity of this point. Therefore it depends strongly on the parameters of


Fig. 7.5. Semiclassical analysis of the population distributions. (a) The energy dependence $E(n)$ increases at $n \rightarrow \infty$ slower, as compared to the coupling $V(n)$ (solid lines). The arrival time $t(n)$ (dashed line), the power index $S_{m}(n)$ of the tunneling probability $e^{-\operatorname{Im} S_{m}(n)}$ (long dash line), and the width $\delta t(n)$ (dash-dot line) are not shown to scale. (b) Population as a function of the level number $n$ and time $t$ for $V(n)=0.42+0.53 n^{0.9}$ and $E(n)=n^{1 / 2}$ as an example.
the last, classically available state, and sometimes it is necessary to replace the integrals by sums at the integer points $x=n$ near the turning points. Another important feature of the wavepacket duration follows from the fact that the integral (7.86) converges at the upper limit if the dependence $V(n)$ rises faster than $n^{1 / 3}$ in the classically accessible region $n \rightarrow \infty$. Therefore $\delta t$ does not depend on the level number for large $n$.

We also note that if the combination $4 V(n)^{2}-E(n)^{2}$ is positive in the limit of large $n$ and tends to infinity faster than $n^{2}$, the integral (7.84) also converges as $n \rightarrow \infty$, and therefore the center of the classical wavepacket reaches infinity at a finite time. In this regime, a level population with a Gaussian time-dependence does not imply a Gaussian distribution over the levels at a given time. On the contrary, this distribution has an essentially non-Gaussian character, being strongly asymmetric, since the higher the number $n$ is, the higher is the local velocity of the populations. This asymmetry can already be seen for $V(x) \sim x^{a}$ for $1 / 2<a<1$, as shown in Fig. 7.5(b).

### 7.3.3 The Tunneling Probability

The wavepacket cannot propagate in the classically forbidden region where $|E(n)-\varepsilon|>|2 V(n)|$. The integrand of (7.79) assumes the imaginary values and the integral

$$
\begin{equation*}
\operatorname{Im}\left(S_{m}\right)=i \int_{n_{b}}^{n_{a}} \arccos \frac{\varepsilon-E(x)}{2 V(x)} d x \tag{7.87}
\end{equation*}
$$

between the points $n_{a, b}$ where $\left|E\left(n_{a, b}\right)-\varepsilon\right|=\left|2 V\left(n_{a, b}\right)\right|$, gives the probability amplitude $e^{-\operatorname{Im}\left(S_{m}\right)}$ for tunneling through the classically forbidden region.

A more instructing form of this expression can be gained if we replace the arccosine function by the corresponding logarithm and the integral by a sum

$$
\begin{align*}
e^{-\operatorname{Im}\left(S_{m}\right)} & =\exp \left[-i \int_{n_{b}}^{n_{a}} \arccos \frac{\varepsilon-E(x)}{2 V(x)} d x\right] \\
& =\exp \left\{\sum_{n=n_{b}}^{n_{a}} \ln \left[\left|\frac{\varepsilon-E(n)}{2 V(n)}\right|-\sqrt{\left(\frac{\varepsilon-E(n)}{2 V(n)}\right)^{2}-1}\right]\right\} \\
& =\prod_{n=n_{b}}^{n_{a}}\left|\frac{|\varepsilon-E(x)|-\sqrt{(\varepsilon-E(x))^{2}-4 V(x)^{2}}}{2 V(x)}\right| \tag{7.88}
\end{align*}
$$

where $n_{a, b}$ now represents the first and the last integer points in the classically forbidden region. Expansion over the small ratio $2 V(x) /(\varepsilon-E(x))$ gives the result

$$
\begin{equation*}
e^{-\operatorname{Im}\left(S_{m}\right)}=\prod_{n=n_{b}}^{n_{a}}\left|\frac{V(n)}{\varepsilon-E(n)}\right| \tag{7.89}
\end{equation*}
$$

which coincides with the first non-vanishing order of perturbation theory. Indeed, as we have seen in (3.17), the amplitude $\psi_{1}$ of a level $|1\rangle$ strongly detuned by an energy $\Delta=\varepsilon-E(1)$ from the energy $\varepsilon$ of the initially populated state $|0\rangle$ amounts at most to a fraction $V(1) /(\varepsilon-E(1))$ of the amplitude $\psi_{0}$. The level $|1\rangle$ serves, in turn, as a starting state for the subsequent set of two levels $|1\rangle$ and $|2\rangle$. Therefore, by analogy, $\psi_{2} \sim V(1) V(2) /(\varepsilon-E(1))(\varepsilon-E(2))$, and repeating the same reasonings for all strongly detuned levels of the classically forbidden region, we arrive at (7.89).

### 7.3.4 Applicability of the WKB Approximation

We note that, contrary to the WKB analysis of the Schrödinger equation with a quadratic dependence of the kinetic energy on the momentum, the zero order WKB approximation cannot be easily improved for the discrete level systems considered. The main reason for this is that a consistent equation for the corrections to the mechanical action (7.79) has the same structure as the original equation (7.76) and contains all-order derivatives. However, from WKB analysis of the regular Schrödinger equation, we know that this correction is important mainly near turning points where the integrals for action correction diverge. The turning points seldom coincide with the integer points, important for the ensembles of discrete levels. Therefore, cut-off or replacement of integrals diverging near the turning points by the corresponding sums may compensate to some extent for the impossibility of improving the zero-order approximation.

The other limitation of the WKB analysis has to do with the fact that, for the validity of the approach, the duration $\delta t(n)(7.86)$ of the wavepacket
at a given state $n$ should be much shorter compared to the arrival time $t(n)$ (7.84) to this level. This is not the case when the coupling increases faster than $\sim n$, when both quantities assume finite values. Neither is it the case for decreasing couplings, when the width increases faster with $n$ than the arrival time. Therefore, the applicability of this approach is restricted to the classically accessible domains where $V(n) \sim n^{a}$ at $n \rightarrow \infty$ with $0<a<1$. It is worth mentioning that in the limit $a=0$, that is for constant coupling, the expression (7.80) is valid and, moreover, it yields solutions of the type (7.10), but the saddle-point approximation of the integrals is incorrect. The population distribution does not have the structure of a Gaussian wavepacket, although the front of this distribution indeed propagates with velocity (7.85).

### 7.4 Relay with disordered parameters

Efficient transport of population from a quantum level to another one, separated by a large interval in the relay-like sequence of coupled states, is possible only if both levels are equally represented in the large number of different eigenfunctions of the system. Formally this statement can be given in terms of the resolvent $1 /(\varepsilon-\widehat{H})(4.10)$ which is another name for Green's operator. In the basis of eigenfunctions $\left|e_{k}\right\rangle$ the Green's operator reads

$$
\begin{equation*}
\frac{1}{\varepsilon-\widehat{H}}=\sum_{k} \frac{\left|e_{k}\right\rangle\left\langle e_{k}\right|}{\varepsilon-\varepsilon_{k}} \tag{7.90}
\end{equation*}
$$

where $\varepsilon_{k}$ denotes energies of the eigenstates, and the amplitude of the transition from the level $n_{0}=0$ to the level $n$ adopts the form

$$
\begin{equation*}
\psi_{n}(t)=\int d \varepsilon e^{-i \varepsilon t} \sum_{k} \frac{\psi_{n}\left(\varepsilon_{k}\right) \psi_{0}^{*}\left(\varepsilon_{k}\right)}{\varepsilon-\varepsilon_{k}} \tag{7.91}
\end{equation*}
$$

where $\psi_{n}\left(E_{k}\right)=\left\langle n \mid e_{k}\right\rangle$ is the amplitude of the level $n$ in the eigenstate $k$. Comparison of (7.91) and (7.41) for $V_{n}=1$ shows a useful correspondence:

$$
\begin{align*}
\psi_{n}\left(\varepsilon_{k}\right) \psi_{0}^{*}\left(\varepsilon_{k}\right) & =\frac{M_{n}^{+}\left(\varepsilon=\varepsilon_{k}\right) M_{0}^{-}\left(\varepsilon=\varepsilon_{k}\right)}{(\partial \operatorname{Det}|\widehat{H}-\varepsilon| / \partial \varepsilon)_{\varepsilon=\varepsilon_{k}}} \\
& =\frac{M_{0}^{+}\left(\varepsilon_{k}\right) M_{0}^{-}\left(\varepsilon_{k}\right)}{(\partial \operatorname{Det}|\widehat{H}-\varepsilon| / \partial \varepsilon)_{\varepsilon=\varepsilon_{k}}} \prod_{1}^{l=n} u_{l}^{+}\left(\varepsilon_{k}\right) \tag{7.92}
\end{align*}
$$

between the eigenfunctions $\psi_{n}\left(\varepsilon_{k}\right)$ and the minors $M_{n}^{+}\left(\varepsilon_{k}\right)$ of the matrix $\widehat{H}-\varepsilon$, where now, in contrast to (7.39), the system is infinite on both sides, and $u_{k}^{+}=M_{k}^{+} / M_{k-1}^{+}$as earlier.

From (7.41) one immediately finds the expression

$$
\begin{align*}
\rho_{n}(t) & =\int_{-\infty-i \nu}^{\infty-i \nu} \frac{d \xi}{2 \pi i} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \varepsilon}{2 \pi i} e^{-i(\varepsilon-\xi) t}(\widehat{H}-\varepsilon)_{k n}^{-1}(\widehat{H}-\xi)_{k n}^{-1} \\
& =\int_{-\infty-i \nu}^{\infty-i \nu} \frac{d \xi}{2 \pi i} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \varepsilon}{2 \pi i} e^{-i(\varepsilon-\xi) t} \frac{M_{k}^{+}(\varepsilon) M_{n}^{-}(\varepsilon)}{\operatorname{Det}|\widehat{H}-\varepsilon|} \frac{M_{k}^{+}(\xi) M_{n}^{-}(\xi)}{\operatorname{Det}|\widehat{H}-\xi|} \tag{7.93}
\end{align*}
$$

for the populations $\rho_{n}(t)=\left|\psi_{n}(t)\right|^{2}$, which with the help of (7.90) can also be given in terms of the eigenfunctions $\psi_{n}\left(\varepsilon_{k}\right)$ of the Hamiltonian

$$
\begin{align*}
\rho_{n}(t) & =\iint d \xi d \varepsilon e^{-i(\varepsilon-\xi) t} \sum_{k, k^{\prime}} \frac{\psi_{n}\left(\varepsilon_{k}\right) \psi_{n}^{*}\left(\varepsilon_{k^{\prime}}\right) \psi_{0}^{*}\left(\varepsilon_{k}\right) \psi_{0}\left(\varepsilon_{k^{\prime}}\right)}{4 \pi^{2}\left(\varepsilon-\varepsilon_{k}\right)\left(\xi-\varepsilon_{k^{\prime}}\right)} \\
& =\int d \zeta e^{-i \zeta t} \sum_{k, k^{\prime}} \frac{\psi_{n}\left(\varepsilon_{k}\right) \psi_{n}^{*}\left(\varepsilon_{k^{\prime}}\right) \psi_{0}^{*}\left(\varepsilon_{k}\right) \psi_{0}\left(\varepsilon_{k^{\prime}}\right)}{2 \pi i\left(\zeta+\varepsilon_{k}-\varepsilon_{k^{\prime}}\right)} \tag{7.94}
\end{align*}
$$

where, as earlier in (3.112), we have introduced the variables $\eta=(\varepsilon+\xi) / 2$ and $\zeta=\varepsilon-\xi$ and have performed integration over $d \eta$. Now one sees, that when the product $\psi_{n}\left(\varepsilon_{k}\right) \psi_{0}^{*}\left(\varepsilon_{k}\right)$ of two amplitudes corresponding to the same energy $\varepsilon_{k}$ of the eigenstate $k$ in the numerator is small, the transition probability is small as well.

It turns out that for the vast majority of the system described by the three-term recurrence equation (7.2), the probability of transition from an initially populated state to another state at some distance $\Delta n$ vanishes at large distances. Indeed, a typical realization of a relay-like system of interacting levels corresponds to the energies $E_{k}$ irregularly changing with the level number $n$. Numerical analysis shows that the eigenfunctions of disordered systems are exponentially localized. This phenomenon is known as Anderson localization. More precisely, for almost all one-dimensional quantum systems, the amplitudes of eigenstates exponentially decrease when $|n| \rightarrow \infty$, whereas all of the exactly solved problems considered earlier in this section are exceptions, resulting from a special underlying group symmetry. In this subsection we give a semiqualitative description of the localization phenomenon, useful for an intuitive understanding of the dynamics of multilevel systems in the opposite limit of total absence of such a symmetry. The main tool of the analysis is the ensemble average which will be performed with a model distribution function of the states energies. The exactly soluble case, known as the Lloyd model, allows one to find the ensemble averaged resolvents for the Cauchy distribution of the level energies. More rigorous analysis based on the resolvents for populations requires a more sophisticated field theory technique whose basic elements are given in the next subsection.

It is worth mentioning that the Hamiltonian $\widehat{H}_{0}=E_{n} \delta_{n, n^{\prime}}$ and the coupling $\widehat{V}=\delta_{n+1, n^{\prime}}+\delta_{n-1, n^{\prime}}$ do not commute. Moreover, for the case of random $E_{n}$ the higer-order commutators remain linearly independent. However, these commutators are not all of the same order of magnitude, since with the
increasing order of commutation $r$, their sizes decrease like $r^{-1 / 2}$, demonstrating the dependence typical of a random process. That is the algebraic reason for the existance of the localization phenomenon considered in this section.

### 7.4.1 Ensemble Averaged Amplitudes and Corresponding Populations

We start with a quantity most convenient for calculation - the mean probability amplitudes $\left\langle\psi_{n}(t)\right\rangle$ for the system to be in a given level $n$. Following P. Lloyd, we assume that in an ensemble of systems (7.2), the level energies $E_{n}$ have independent, Lorentzian (Cauchy) distributions

$$
\begin{equation*}
\widetilde{g}\left(E_{n}\right)=\frac{\gamma / \pi}{\left(E_{n}-\Omega\right)^{2}+\gamma^{2}}, \tag{7.95}
\end{equation*}
$$

where $\left|V_{n}\right|^{2}=1$. This situation is shown in Fig. 7.6(a).
The resolvent $1 /(\varepsilon-\widehat{H})$ of (7.2)

$$
0=\left(E_{n}-\varepsilon\right) \psi_{n}+\psi_{n-1}+\psi_{n+1}
$$

given by the explicit expression (7.41) is a regular function in the lower part of complex planes of all $E_{n}$, since the Fourier variable $\varepsilon$ is shifted to the upper part of the complex plane by a vanishing value $\nu \rightarrow 0$, as shown in Fig. 3.2. Therefore calculation of the average with the distributions (7.95) can simply be done by taking the residuals at the points $E_{n}=\Omega-i \gamma$ which means the replacements $E_{n} \rightarrow \Omega-i \gamma$ in the resolvent. But such an average resolvent corresponds to the Schrödinger equation

$$
\begin{equation*}
i \frac{\partial \psi_{n}}{\partial t}=(\Omega-i \gamma) \psi_{n}+\psi_{n-1}+\psi_{n+1} \tag{7.96}
\end{equation*}
$$

which after the replacement $\psi_{n}(t) \rightarrow \psi_{n}(t) e^{-\gamma t}$ coincides with (7.1) for $E_{n}=$ $\Omega, V_{n}=1$. With the help of the solution (7.11), for (7.96) we immediately find

$$
\begin{equation*}
\psi_{n}(t)=i^{n} e^{-\gamma t-i \Omega t} J_{n}(2|t|) \tag{7.97}
\end{equation*}
$$

The ensemble averaged equation (7.96) coincides with the Schrödinger equation of a cascade of resonant levels, each of which is coupled to an independent uniform continuum, by analogy to the case of a single level coupled to a continuum considered in Sect.3.2.2 as shown in Fig. 7.6(b). Let us calculate the harmonics of the total "fluence"

$$
\begin{equation*}
\Phi_{n}(\varepsilon)=\gamma \int_{0}^{\infty} e^{i \varepsilon t} \psi_{n}(t) d t=i^{n}\left[4+(\gamma+i \Omega-i \varepsilon)^{2}\right]^{-(n+1) / 2} \tag{7.98}
\end{equation*}
$$



Fig. 7.6. Anderson localization in one dimension and its simplistic description by the mean Green's function. (a) A set of levels with the nearest neighbor coupling and disordered level positions given by the distribution function $\tilde{g}\left(E_{n}\right)$ of width $\gamma$ suggested by (7.95). (b) Naive description assuming that the Lorentzian distribution of the level positions is equivalent to level dacay to continua with a rate $\gamma$. (c) Distribution of the populations (7.100) suggested by the ensemble averaged resolvent. (d) Density of states of the disordered chain given by (7.103) for $\gamma=0.2$ (solid line), $\gamma=1$ (dashed line), $\gamma=2.5$ (dash-dot line). For $\gamma=0$ the density of states tends to infinity at the edges of the spectrum $\epsilon \rightarrow \pm 2$.
of the population amplitude to the continua, where in order to perform the integration we have employed the Hankel integral representation of the Bessel function. We can also find the asymptotic populations of the continua

$$
\begin{equation*}
\widetilde{\rho}_{n}=2 \gamma \int_{0}^{\infty}\left|\psi_{n}(t)\right|^{2} d t=\frac{\Gamma\left(n+\frac{1}{2}\right)}{\gamma^{2 n} n!\sqrt{\pi}}{ }_{2} F_{1}\left(n+\frac{1}{2}, n+\frac{1}{2} ; 1+2 n ; \frac{-4}{\gamma^{2}}\right) \tag{7.99}
\end{equation*}
$$

that coincide with $\left|\Phi_{n}(\varepsilon)\right|^{2}$ integrated over the spectrum $\varepsilon$. In order to gain a deeper insight into the population dynamics in Fig. 7.6(c) we depict the sum of the population of the levels and corresponding continua

$$
\begin{equation*}
\rho_{n}(t)=\left|\psi_{n}(t)\right|^{2}+2 \gamma \int_{0}^{t}\left|\psi_{n}(t)\right|^{2} d t \tag{7.100}
\end{equation*}
$$

as a function of time, which shows how the asymptotic distribution (7.99) is attained. One sees that the population initially localized at level $n=0$ spreads over the neighboring states with the course of time, until it reaches a steady state distribution, exponentially decreasing for $|n| \rightarrow \infty$.

### 7.4.2 Ensemble Averaged Spectrum

We can also find the ensemble averaged spectrum density of the eigenstates. To this end we note that the energy spectrum of a quantum system can be found from the resolvent with the help of the expression

$$
\begin{align*}
g(\varepsilon) & =\frac{-1}{\pi} \lim _{\nu \rightarrow 0} \operatorname{Im}\left(\operatorname{Tr} \frac{1}{\varepsilon-\widehat{H}+i \nu}\right) \\
& =\frac{-1}{\pi} \lim _{\nu \rightarrow 0} \operatorname{Im}\left(\sum_{k} \frac{\left\langle e_{k} \mid e_{k}\right\rangle}{\varepsilon-\varepsilon_{k}+i \nu}\right)=\sum_{k} \delta\left(\varepsilon-\varepsilon_{k}\right) . \tag{7.101}
\end{align*}
$$

The simplest way to calculate the average spectrum of a disordered chain is to express the Hamiltonian with the help of the momentum operator $p=-i \frac{\partial}{\partial n}$ which we have already employed when considering the classical limit (7.78)

$$
\begin{equation*}
\widehat{H}=E_{n}+2 \cos \widehat{p}, \tag{7.102}
\end{equation*}
$$

and obtain after integration over the energies $E_{n}$ with the distributions (7.95)

$$
\begin{align*}
g(\varepsilon) & =\frac{-1}{\pi} \operatorname{Im}\left(\operatorname{Tr} \frac{1}{\varepsilon-\Omega+i \gamma-2 \cos \widehat{p}}\right) \\
& =\frac{-1}{\pi} \operatorname{Im} \int_{-\pi}^{\pi} \frac{d p}{\varepsilon-\Omega+i \gamma-2 \cos p} \\
& =\operatorname{Re} \frac{2}{\sqrt{4-(\varepsilon-\Omega+i \gamma)^{2}}} \tag{7.103}
\end{align*}
$$

where in the last line we have taken advantage of the invariance of the trace operator and have employed the momentum representation. One sees that this spectrum in Fig. 7.6(d). It is given by the inverse of the combination $\sqrt{4-(\varepsilon-\Omega+i \gamma)^{2}}$, the same as that which determines the amplitude fluence distribution (7.98).

### 7.4.3 Distribution of the Amplitude Ratios

Strictly speaking, the population distribution (7.100) calculated with the help of the ensemble averaged amplitudes has nothing to do with the ensemble averaged populations in the relay-like system with random positions of the
levels. It just gives an "eye guiding line" and some parameters that might be of use for a consistent consideration, which is technically much more involved. Let us consider this phenomenon more carefully with the help of a different technique. From (7.2)

$$
0=\left(E_{n}-\varepsilon\right) \psi_{n}(\varepsilon)+\psi_{n-1}(\varepsilon)+\psi_{n+1}(\varepsilon)
$$

and from the distribution of the level positions (7.95) we find distribution functions $g_{n}^{ \pm}\left(u_{n}^{ \pm}\right)$of the ratios

$$
\begin{equation*}
u_{n+1}^{+}(\varepsilon)=\psi_{n+1}(\varepsilon) / \psi_{n}(\varepsilon) \tag{7.104}
\end{equation*}
$$

and

$$
\begin{equation*}
u_{n-1}^{-}(\varepsilon)=\psi_{n-1}(\varepsilon) / \psi_{n}(\varepsilon) \tag{7.105}
\end{equation*}
$$

that satisfy (7.46)

$$
\begin{aligned}
& u_{n+1}^{+}=\left(\varepsilon-E_{n}\right)-\frac{1}{u_{n}^{+}} \\
& u_{n-1}^{-}=\left(\varepsilon-E_{n}\right)-\frac{1}{u_{n}^{-}} .
\end{aligned}
$$

Note that we can also have in mind the ratios of minors of (7.92), as in the original form of (7.46).

We assume the distribution (7.95) for the level energies $\widetilde{g}\left(E_{n}\right)$, and ask for the distributions $g_{n}^{ \pm}\left(u_{n}^{ \pm}\right)$of the ratios $u_{n}^{ \pm}$. Due to the physical meaning of the probability distribution, the functions $g_{n}^{ \pm}(x)$ have to be real and positive at any complex value of the argument. For the sake of simplicity, let us consider only real $x$. With the help of (7.46) we arrive at the recurrence equation

$$
\begin{align*}
g_{n+1}^{+}\left(u_{n+1}^{+}\right) & =\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \delta\left[u_{n+1}^{+}-\varepsilon+E_{n}+\frac{1}{u_{n}^{+}}\right] g_{n}^{+}\left(u_{n}^{+}\right) \widetilde{g}\left(E_{n}\right) d E_{n} d u_{n}^{+} \\
g_{n+1}^{+}(x) & =\int_{-\infty}^{\infty} \frac{\gamma / \pi}{(x-\varepsilon+1 / y+\Omega)^{2}+\gamma^{2}} g_{n}^{+}(y) d y \tag{7.106}
\end{align*}
$$

where in the last line we have replaced the integration variable $u_{n}^{+}$by $y$ and $u_{n+1}^{+}$by $x$. A similar equation holds for $g_{n}^{-}\left(u_{n}^{-}\right)$.

The functional equation (7.106) may have many different complicated solutions. We will consider the simplest solution which has the structure

$$
\begin{equation*}
g_{n}^{+}(y)=\frac{\gamma_{n} / \pi}{\left(y-\omega_{n}\right)^{2}+\gamma_{n}^{2}} . \tag{7.107}
\end{equation*}
$$

In other words, we assume that distribution functions for different $n$ have the same functional structure, and differ only by the positions $\omega_{n}$ and width $\gamma_{n}$ of the single spike. Substitution of (7.107) into (7.106), followed by the integration around the poles in the upper part of the complex plane, yields

$$
\begin{equation*}
g_{n+1}^{ \pm}(y)=\frac{\gamma_{n+1} / \pi}{\left(y-\omega_{n+1}\right)^{2}+\gamma_{n+1}^{2}} \tag{7.108}
\end{equation*}
$$

where

$$
\begin{align*}
\gamma_{n+1} & =\frac{\gamma_{n}}{\omega_{n}^{2}+\gamma_{n}^{2}}+\gamma \\
\omega_{n+1} & =\frac{-\omega_{n}}{\omega_{n}^{2}+\gamma_{n}^{2}}+\Omega+\varepsilon \tag{7.109}
\end{align*}
$$

The last equations can be written in a more compact form

$$
\begin{equation*}
z_{n+1}=\frac{-1}{z_{n}}+Z \tag{7.110}
\end{equation*}
$$

when we introduce complex numbers $z_{n}=\omega_{n}+i \gamma_{n}$ and $Z=\Omega-\varepsilon+i \gamma$. For $g_{n}^{-}(y)$ one obtains a similar relation

$$
\begin{equation*}
z_{n-1}=\frac{-1}{z_{n}}+Z \tag{7.111}
\end{equation*}
$$



Fig. 7.7. Transformation of the complex plane given by (7.110). A complex number $z_{n}$ (shown by the solid vector) after taking the inverse value $1 / z_{n}$, change of sign (dotted vectors), and after displacement by the complex number $Z$ (dashed vector), becomes a complex number $z_{n+1}$. The stable stationary point of this transformation (center of the small solid circle) is located outside of the unit circle, whereas the inverse value corresponds to the unstable stationary point (center of the dotted circle).

Equations (7.110), (7.111) map the complex plane $z$ to itself. There is a clear geometric meaning, as shown in Fig.7.7. The transformation $-1 / z$
corresponds to inversion of a complex point $z$ with respect to the unit circle, followed by the sign inversion of the real part, whereas addition of the complex number $Z$ means displacement of this point. This mapping has two stationary points, given by the solution of the algebraic equation

$$
\begin{equation*}
z=\frac{-1}{z}+Z \tag{7.112}
\end{equation*}
$$

that corresponds to the situation where displacement of the point as a result of inversion and reflection with respect to the imaginary axis exactly equals $-Z$. The solution of (7.16) suggests

$$
\begin{equation*}
z_{ \pm}=\frac{Z}{2} \pm \sqrt{\frac{Z}{4}^{2}-1} \tag{7.113}
\end{equation*}
$$

The stationary point $z_{+}$with $\left|z_{+}\right|>1$ is stable for (7.110), whereas $\left|z_{-}\right|<1$ corresponds to an unstable stationary point $z_{-}$. This follows from the fact that inversion to the inner part of a circle is a compressing mapping. And for (7.111), the situation is the opposite. Note that the indices $\pm$ of $z_{ \pm}$do not always correspond to the sign of the square root in (7.113), which for complex $Z$ may depend on the choise of the square root branch, and hence the condition $\left|z_{+}\right|>1$ should be checked.

The presence of stationary points means that, after several iterations, any initial Lorentzian distribution of the ratios (7.104), (7.105) results in a Lorentzian distribution of width $\operatorname{Im} z_{ \pm}$around the points $\operatorname{Re} z_{ \pm}$. Moreover, the linearity of (7.106) implies that any distribution that can be represented as a sum of Lorentzians adopts this form. Therefore, any analytical distribution that has an arbitrary number of simple poles, placed symmetrically with respect to the real axis, becomes Lorentzian

$$
\begin{equation*}
g_{n}^{ \pm}\left(u_{n}^{ \pm}\right)=\frac{\operatorname{Im} z_{+} / \pi}{\left(u_{n}^{ \pm}-\operatorname{Re} z_{+}\right)^{2}+\operatorname{Im} z_{+}^{2}} . \tag{7.114}
\end{equation*}
$$

With the help of the distribution (7.107) we can find a mean value which will play an important role in our further analysis. We calculate the mean $\log \left(u_{n}^{ \pm}\right)$and concentrate on its real part. Straightforward integration yields

$$
\begin{align*}
\langle\ln | u_{n}^{ \pm}(\varepsilon)| \rangle & =\frac{1}{\pi} \int_{-\infty}^{\infty} d u_{n}^{ \pm} \frac{\ln \left|u_{n}^{ \pm}(\varepsilon)\right| \operatorname{Im} z_{+}}{\left(u_{n}^{ \pm}-\operatorname{Re} z_{+}\right)^{2}+\left(\operatorname{Im} z_{+}\right)^{2}} \\
& =\ln \left|z_{+}(\varepsilon)\right|=\ln \left|\frac{Z}{2} \pm \sqrt{\frac{Z^{2}}{4}-1}\right| \tag{7.115}
\end{align*}
$$

where the sign of the square root should be chosen such that $\left|z_{+}(\varepsilon)\right|>1$. Substitution of $Z=2 \cosh X$ yields another form of this expression

$$
\begin{align*}
\ln \left|z_{+}(\varepsilon)\right| & =|\operatorname{Re} X|=\left|\operatorname{Re}\left(\operatorname{arc} \cosh \frac{Z}{2}\right)\right| \\
& =\operatorname{arccosh}\left[\frac{1}{2}\left|1-\frac{Z}{2}\right|+\frac{1}{2}\left|1+\frac{Z}{2}\right|\right] . \tag{7.116}
\end{align*}
$$

Note that the mean values (7.116) of both $\log u_{n}^{+}$and $\log u_{n}^{-}$are positive, although the product $u_{n}^{+} u_{n}^{-}$is equal to unity for any specific realization of the system, as follows from the definitions (7.104)-(7.105). There is no contradiction in this fact, since (7.2) for a generic $\varepsilon$ has two solutions - one exponentially increasing for $n \rightarrow \infty$ and the other exponentially increasing for $n \rightarrow-\infty$. Only when $\varepsilon$ equals an eigenvalues $\varepsilon_{k}$ of a specific realization of the random system, do these two solutions match at some point $n$, but only for this specific realization. At a fixed $\varepsilon$, this does not happen for a typical realization in the ensemble, and that is why $\langle\ln | u_{n}^{+}| \rangle+\langle\ln | u_{n}^{-}| \rangle \neq\langle\ln | u_{n}^{+} u_{n}^{-}| \rangle$ but tend to the largest stationary values.

Equation (7.115) gives the idea of the typical behavior of the eigenstate amplitudes $\psi_{n}\left(\varepsilon_{k}\right)$ as a function of the level number $n$. These functions are typically growing with the growth rate $\log \left|z_{+}\left(\varepsilon_{k}\right)\right|$ when the number $n$ comes from $\pm \infty$ to a certain finite number $n_{\max }\left(\varepsilon_{k}\right)$ where for a particular realization of the random system, these functions match the three-term recurrence relation. This exponential behavior is typical of all eigenfunctions and of any particular realization of the Hamiltonian, whereas the precise value of the energy $\varepsilon_{k}$ and specific position $n_{\max }\left(\varepsilon_{k}\right)$ are sensitive to the particular realization.


Fig. 7.8. Anderson localization in 1D systems. (a) Moduli of eigenvectors (dots, not to scale) as functions of the level number $n$ are depicted for three values of eigenenergies $\epsilon_{k}$ for $\gamma=0.8$. The values of $\epsilon_{k}$ serve as reference levels of the corresponding plots. The mean state density $g\left(\epsilon_{k}\right)$ (dashed line) of (7.103) and mean ratios of neighboring amplitudes $z_{+}\left(\epsilon_{k}\right)$ (long dash line) of (7.116) are also shown, but not to scale. (b) "Naive" distribution of the populations (7.100) (solid line) compared to the population distribution (7.124) (dash-dot line) for $\gamma=0.15$.

In Fig. 7.8(a) we show three typical examples of the eigenvector amplitudes, corresponding to different energies of the eigenstates. One clearly sees the exponential decrease of the amplitudes. The ensemble averaged density of states and the mean ratios of the neighboring amplitudes are also shown as a function of energy.

### 7.4.4 Distribution of the Populations for Long Times

Let us now consider the population distribution (7.94) which has the form of a product of two mutually complex conjugate probability amplitudes. After the ensemble average, the sums over the eigenenergies $\varepsilon_{k}$ and $\varepsilon_{k^{\prime}}$ can be replaced by integrals over their distributions $g\left(\varepsilon_{k}, \varepsilon_{k^{\prime}}\right)$ which can always be represented as a $\operatorname{sum} g\left(\varepsilon_{k}, \varepsilon_{k^{\prime}}\right)=g\left(\varepsilon_{k}\right) g\left(\varepsilon_{k^{\prime}}\right)+g_{2}\left(\varepsilon_{k}, \varepsilon_{k^{\prime}}\right)$ of the independent binary distribution $g\left(\varepsilon_{k}\right) g\left(\varepsilon_{k^{\prime}}\right)$, given by (7.103), and a correlation term $g_{2}\left(\varepsilon_{k}, \varepsilon_{k^{\prime}}\right)$. This yields

$$
\begin{align*}
\rho_{n}(t)= & \iint e^{-i \varepsilon t} g\left(\varepsilon_{k}\right) \frac{\psi_{n}\left(\varepsilon_{k}\right) \psi_{0}^{*}\left(\varepsilon_{k}\right)}{2 \pi i\left(\varepsilon-\varepsilon_{k}\right)} d \varepsilon_{k} d \varepsilon \\
& \times \iint e^{i \xi t} g\left(\varepsilon_{k^{\prime}}\right) \frac{\psi_{n}^{*}\left(\varepsilon_{k^{\prime}}\right) \psi_{0}\left(\varepsilon_{k^{\prime}}\right)}{-2 \pi i\left(\xi-\varepsilon_{k^{\prime}}\right)} d \varepsilon_{k^{\prime}} d \xi \\
& +\iiint g_{2}\left(\varepsilon_{k}, \varepsilon_{k^{\prime}}\right) e^{-i \zeta t} \frac{\psi_{n}\left(\varepsilon_{k}\right) \psi_{n}^{*}\left(\varepsilon_{k^{\prime}}\right) \psi_{0}^{*}\left(\varepsilon_{k}\right) \psi_{0}\left(\varepsilon_{k^{\prime}}\right)}{2 \pi i\left(\zeta+\varepsilon_{k}-\varepsilon_{k^{\prime}}\right)} d \zeta d \varepsilon_{k} d \varepsilon_{k^{\prime}} \tag{7.117}
\end{align*}
$$

where in the last term, the variables $\eta=(\varepsilon+\xi) / 2$ and $\zeta=\varepsilon-\xi$ of (7.94) have been employed.

The first term of (7.117) has the clear physical meaning of the product of two ensemble averaged Green operators that result in the term $\left|\psi_{n}(t)\right|^{2}$ of (7.100), whereas the last term allows for the correction $\widetilde{\rho}_{n}(t)$ which in (7.100) has been replaced by the heuristic term $2 \gamma \int_{0}^{t}\left|\psi_{n}(t)\right|^{2} d t$. We concentrate on this last term, since it is responsible for the stationary population distribution which is attained in the long-time limit. With the help of (7.104) and (7.105) we write this term in the form

$$
\begin{equation*}
\widetilde{\rho}_{n}(t)=\iiint \frac{\left|\psi_{0}\left(\varepsilon_{k}\right)\right|^{2}\left|\psi_{0}\left(\varepsilon_{k^{\prime}}\right)\right|^{2}}{2 \pi i\left(\zeta+\varepsilon_{k}-\varepsilon_{k^{\prime}}\right)} \frac{e^{-i \zeta t} g_{2}\left(\varepsilon_{k}, \varepsilon_{k^{\prime}}\right) d \zeta d \varepsilon_{k} d \varepsilon_{k^{\prime}}}{\prod_{m=1}^{n-1} u_{m}^{-}\left(\varepsilon_{k}\right)\left(u_{m}^{-}\left(\varepsilon_{k^{\prime}}\right)\right)^{*}} \tag{7.118}
\end{equation*}
$$

which shows that, as $t \rightarrow \infty$, the dominating time-independent contribution comes from $k=k^{\prime}$. This yields

$$
\begin{equation*}
\left.\tilde{\rho}_{n}\right|_{t \rightarrow \infty}=\int\left|\psi_{0}\left(\varepsilon_{k}\right)\right|^{4} g_{2}\left(\varepsilon_{k}, \varepsilon_{k}\right) \prod_{m=1}^{n-1}\left|u_{m}^{-}\left(\varepsilon_{k}\right)\right|^{-2} d \varepsilon_{k} \tag{7.119}
\end{equation*}
$$

whereas the time-dependent part reads

$$
\begin{equation*}
\tilde{\rho}_{n}(t)-\left.\tilde{\rho}_{n}\right|_{t \rightarrow \infty}=\int g_{n}^{(2)}(\zeta) e^{-i \zeta t} d \zeta \tag{7.120}
\end{equation*}
$$

where the function

$$
\begin{equation*}
g_{n}^{(2)}(\zeta)=\int d x \frac{\left|\psi_{0}\left(x+\frac{\zeta}{2}\right)\right|^{2}\left|\psi_{0}\left(x-\frac{\zeta}{2}\right)\right|^{2} g_{2}\left(x+\frac{\zeta}{2}, x-\frac{\zeta}{2}\right)}{\prod_{m=1}^{n-1} u_{m}^{-}\left(x+\frac{\zeta}{2}\right)\left(u_{m}^{-}\left(x-\frac{\zeta}{2}\right)\right)^{*}} \tag{7.121}
\end{equation*}
$$

has the physical meaning of the correlation of the energy distance between eigenstates, weighted by the correlations of the probability amplitudes to be in the initial and the final levels corresponding to these states. Note that the joint distribution $g_{2}\left(x, x^{\prime}\right)$ of two eigenvalues is discontinuous at $x=x^{\prime}$, assuming at this point the dependence (7.103) of a single eigenvalue distribution $g(x)$. Therefore, from the integral (7.120) we have to exclude the point $\zeta=0$ which is included in the part (7.119) with $g_{2}\left(\varepsilon_{k}, \varepsilon_{k}\right)=g\left(\varepsilon_{k}\right) N$, where $N$ is the total number of levels.

Consider first the stationary distribution (7.119). According to (7.115)

$$
\begin{equation*}
\prod_{m=1}^{n-1}\left|u_{m}^{-}\left(\varepsilon_{k}\right)\right|^{-2}=e^{-2 \sum_{m=1}^{n-1} \ln \left|u_{m}^{-}\left(\varepsilon_{k}\right)\right|} \simeq\left|z_{+}\left(\varepsilon_{k}\right)\right|^{-2|n|} \tag{7.122}
\end{equation*}
$$

since the sum of many randomly distributed logarithms can be replaced by their mean value multiplied by the total number of terms. Therefore

$$
\begin{equation*}
\left.\tilde{\rho}_{n}\right|_{t \rightarrow \infty} \simeq \int\left|\psi_{0}(x)\right|^{4} N g(x)\left|z_{+}(x)\right|^{-2|n|} d x \tag{7.123}
\end{equation*}
$$

where the procedure of taking the average of $\left|\psi_{0}(x)\right|^{4}$ absorbs the deviations of sums of logarithms in (7.122) from their mean values $\left.\left.\langle | \psi_{0}(x)\right|^{4}\right\rangle=$ $\left.\left.\langle | \psi_{0}(x)\right|^{4} \prod_{m=1}^{n-1}\left|z_{+}(x) / u_{m}^{-}\left(\varepsilon_{k}\right)\right|^{2}\right\rangle$ and should be done carefully. To this end, let us take into account the normalization $\left.\sum_{n} \widetilde{\rho}_{n}\right|_{t \rightarrow \infty}=1$ of the population distribution and establish with the help of this condition a relation between the average $\left.\left.\langle | \psi_{0}(x)\right|^{2}\right\rangle$ and $\left|z_{+}(x)\right|$. If $\left.\left.\langle | \psi_{0}(x)\right|^{2}\right\rangle$ entering the product $\left.\left.\langle | \psi_{0}(x)\right|^{2}\right\rangle\left|z_{+}(x)\right|^{-2 n}$ absorbs the deviations $\sum_{m=1}^{n-1} \ln \left|u_{m}^{-}\left(\varepsilon_{k}\right) / z_{+}(x)\right|^{2}$, then the product equals the mean population $\left.\widetilde{\rho}_{n}=\left.\langle | \psi_{0}(x)\right|^{2} \prod_{m=1}^{n-1}\left|u_{m}^{-}\left(\varepsilon_{k}\right)\right|^{-2}\right\rangle$ of the level $n$ corresponding to a given eigenenergy $x$. Therefore the sum $\left.\left.\sum_{n}\langle | \psi_{0}(x)\right|^{2}\right\rangle\left|z_{+}(x)\right|^{-2 n}$ equals 1 , whereas the average of the remaining factor $\left|\psi_{0}(x)\right|^{2}$ in (7.123) is uniform and equals the inverse total number $N$ of levels in the system, thus normalizing the distribution function of the eigenvalues $N g(x)$ to unity. This yields the stationary ensemble averaged population distribution

$$
\begin{equation*}
\left.\widetilde{\rho}_{n}\right|_{t \rightarrow \infty} \simeq \int \frac{g(x)\left|z_{+}(x)\right|^{-2|n|}}{\sum_{k=-\infty}^{\infty}\left|z_{+}(x)\right|^{-2|k|}} d x=\int \frac{1-\left|z_{+}(x)\right|^{-2}}{1+\left|z_{+}(x)\right|^{-2}} g(x)\left|z_{+}(x)\right|^{-2|n|} d x \tag{7.124}
\end{equation*}
$$

which is systematically slightly broader than our first "naive" guess (7.99), as if the distribution (7.124) results from a dynamical process which takes into account partial recurrences from the continua shown in Fig. 7.6(b). In Fig. 7.8(b) we depict both the "naive" guess and the distribution (7.124) calculated with the help of (7.103) and (7.113).

### 7.4.5 Dynamics of the Asymptotic Populations

We consider next the dynamics of populations at long $t$. According to (7.120), the population approaches its stationary distribution (7.124) following a time dependence given by the Fourier transformation of the correlator (7.121). It can also be given in the form

$$
\begin{equation*}
g_{n}^{(2)}(\zeta) \simeq \int d x \frac{\left|\psi_{0}\left(x+\frac{\zeta}{2}\right)\right|^{2}\left|\psi_{0}\left(x-\frac{\zeta}{2}\right)\right|^{2} g_{2}\left(x+\frac{\zeta}{2}, x-\frac{\zeta}{2}\right)}{\left|z_{+}(x)\right|^{2 n-2}} \tag{7.125}
\end{equation*}
$$

where, as earlier, we have replaced the amplitude ratios by their mean values, and have ignored the difference between $u_{m}^{-}\left(x+\frac{\zeta}{2}\right)$ and $u_{m}^{-}\left(x-\frac{\zeta}{2}\right)$ in the limit $\zeta \rightarrow 0$. Indeed, in the limit of long times, the most important domain corresponds to $\zeta \sim 0$. Therefore the main contribution arises from close eigenenergies in the system. The integrand (7.125) can be interpreted as the spectral correlation $g_{2}$ weighted by the joint probability

$$
w=\left|\psi_{0}\left(x+\frac{\zeta}{2}\right)\right|^{2}\left|\psi_{n}\left(x-\frac{\zeta}{2}\right)\right|^{2}
$$

to be at the levels $|0\rangle$ and $|n\rangle$ corresponding to energy eigenstates $\left|x+\frac{\zeta}{2}\right\rangle$ and $\left|x-\frac{\zeta}{2}\right\rangle$, respectively. In a large, disordered system, the probability $g_{2}\left(x+\frac{\zeta}{2}, x-\frac{\zeta}{2}\right)$ to have two close eigenvalues, $x+\zeta / 2$ and $x-\zeta / 2$, is high. However, typically these states are localized in different domains of $n$, and therefore the weight factor $w$ is exponentially small for most of $x$. On the other hand, the eigenstates localized in the same domain of $n$ and having large $w$, typically correspond to distinct eigenenergies. The probability of finding close eigenvalues increases linearly with the increasing size of the considered domain of $n$, whereas their contribution to the integral (7.125) drops exponentially because of the decreasing weight factors $w$. The interplay between these two tendencies forms the asymptotic spectral dependence $g_{n}^{(2)}(\zeta)$ as $\zeta \rightarrow 0$.

Let us discuss a possible form of this asymptotic dependence. On the axis of $n$, consider a long interval of $N$ successive levels symmetrically placed around the point $n=0$. The density of states in the energy domain $x$ amounts to $N g(x)$ and the mean spacing between the neighboring energy eigenstates localized in this domain is $\delta \sim 1 / N g(x)$. Smaller spacings in larger intervals result from the eigenstates localized in the domain with $n>N / 2$. Therefore, in order to have the mean level spacing smaller than $\zeta$, one has to consider levels at long distance with $n>1 / 2 g(x)|\zeta|$. But the overall contribution of these levels to the integral (7.125) is exponentially small, and can be estimated as a double sum

$$
\begin{align*}
C(x, \zeta) & =\sum_{k, k^{\prime}>1 / 2 g(x) \zeta} A^{2} e^{-\ln \left|z_{+}(x)\right|\left(k+k^{\prime}\right)}=\frac{A^{2} e^{-2 \ln \left|z_{+}(x)\right| / g(x) \zeta}}{\left(1-\left|z_{+}(x)\right|^{-2}\right)^{2}} \\
& =\left(1+\left|z_{+}(x)\right|^{-2}\right)^{2} \exp \left[-\frac{\ln \left|z_{+}(x)\right|}{g(x)|\zeta|}\right] \tag{7.126}
\end{align*}
$$

where the normalization $A=\left(1-\left|z_{+}(x)\right|^{-2}\right) /\left(1+\left|z_{+}(x)\right|^{-2}\right)$, similar to that in (7.124), has been taken into account.

The derivative $\partial / \partial \zeta$ of the total contribution $C(x, \zeta)$ of all resonances closer than $\zeta$ gives the weighted correlation $w g_{n}^{(2)}(x, \zeta)$ for a given $x$ and hence

$$
\begin{equation*}
g_{n}^{(2)}(\zeta) \simeq \int d x \frac{\left(1+\left|z_{+}(x)\right|^{-2}\right)^{2}}{\left|z_{+}(x)\right|^{2 n+2}} \frac{\ln \left|z_{+}(x)\right|}{g(x) \zeta^{2}} \exp \left[-\frac{\ln \left|z_{+}(x)\right|}{g(x)|\zeta|}\right] \tag{7.127}
\end{equation*}
$$

Therefore, for the time dependence (7.120), by performing the inverse Fourier transformation with the help of (7.103) and (7.113) one obtains

$$
\begin{equation*}
\tilde{\rho}_{n}(t)-\left.\widetilde{\rho}_{n}\right|_{t \rightarrow \infty}=\int d x \frac{\left(1+\left|z_{+}(x)\right|^{-2}\right)^{2}}{\left|z_{+}(x)\right|^{2 n+2}} \frac{\ln \left|z_{+}(x)\right|}{g(x)} \int_{0}^{\infty} \exp \left[-\frac{\ln \left|z_{+}(x)\right|}{g(x)|\zeta|}\right] \cos (\zeta t) \frac{d \zeta}{\zeta^{2}} \tag{7.128}
\end{equation*}
$$

and to be consistent within the order of approximation, after exact representation of the last integral

$$
\begin{equation*}
\tilde{\rho}_{n}(t)-\left.\tilde{\rho}_{n}\right|_{t \rightarrow \infty}=\int d x \frac{\left(1+\left|z_{+}(x)\right|^{-2}\right)^{2}}{\left|z_{+}(x)\right|^{2 n+2}} 2 \sqrt{t \frac{\log \left|z_{+}(x)\right|}{g(x)}} \operatorname{Re} K_{1}\left[2 \sqrt{i t \frac{\log \left|z_{+}(x)\right|}{g(x)}}\right] \tag{7.129}
\end{equation*}
$$

in terms of the Bessel function $K_{1}(x)$, one has to take the asymptotic of long times which implies $t / n g(x=0) \gg 1$.

Equation (7.129) yields the dying oscillations at a frequency decreasing in time as $t^{-1 / 2}$ and of an amplitude changing as $t^{-1 / 4} \exp [-$ const $\sqrt{t}]$. In order to obtain the last dependence we have taken into account that the main contribution to the integral (7.129) over $x$ at $t \rightarrow \infty$ comes from the band
center, that is from the domain $x \sim 0$ where $\left|z_{+}(x)\right|-\left|z_{+}(0)\right| \propto|x|$, as one sees in Fig. 7.8(a). The localization length $L=1 / 2 \log \left|z_{+}(0)\right|$ assumes in this spectral region the longest possible value. In Sect. 3.4, we have already seen that strongly non-homogeneous bands may result in non-exponential decay laws. The $e^{-\sqrt{t}}$ dependence obtained here indicates once again the fact that the evolution of a disordered system may differ considerably from the standard exponential relaxation, typical of continua. The common feature of such systems is an infinitely dense but not continuous spectrum, where the eigenfunctions and related matrix elements of physical operators are discontinuous at every point and may change considerably, by orders of magnitudes, with infinitesimal changes of the energy.

### 7.5 Field Theory Method for Disordered Systems

Considering in Sect.7.2.1, the Anderson localization for a one-dimensional disordered relay of levels, we have employed the Cauchy distribution (7.95) of the level energies. This particular type of ensemble has enabled us to determine explicitly the distribution (7.106) of the neighboring amplitude ratios, which has been further employed for the analysis of the population dynamics. The analysis, however, relied on qualitative reasonings rather than on straightforward calculations of the physical quantities of interest. Though this approach helps to develop an intuitive vision of the problem, one needs another, based on straightforward calculations that would leave no doubts about the results obtained. It turns out that this goal is achieved with the help of a Gaussian distribution of the level energies

$$
\begin{equation*}
\widetilde{g}\left(E_{n}\right)=\frac{1}{\sqrt{\pi} \gamma} \exp \frac{-\left(E_{n}-\Omega\right)^{2}}{\gamma^{2}} \tag{7.130}
\end{equation*}
$$

The required technique for this calculation is well known in field theory. In this subsection we discuss it, trying to keep the narration close (as much as possible) to the manner of presentation of the rest of the text.

### 7.5.1 Tunneling Transparency and Classical Bosonic Fields

We concentrate first on the simplest problem, and consider the tunneling transparency of a disordered chain. Equation (7.41) yields the probability

$$
\begin{equation*}
\rho_{N}(t)=\int_{-\infty-i \nu}^{\infty-i \nu} \frac{d \xi}{2 \pi i} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \varepsilon}{2 \pi i} \frac{\exp [-i t(\varepsilon-\xi)]}{\operatorname{Det}|\widehat{H}-\varepsilon| \operatorname{Det}|\widehat{H}-\xi|} \tag{7.131}
\end{equation*}
$$

to find a particle at the last level $N$ of the chain at a time $t$, provided that at $t=0$ it has been at the first level. The Hamiltonian

$$
\widehat{H}=\left\|\begin{array}{ccccccc}
E_{1} & 1 & 0 & . & 0 & 0 & 0  \tag{7.132}\\
1 & E_{2} & 1 & . & 0 & 0 & 0 \\
0 & 1 & E_{3} & . & 0 & 0 & 0 \\
. & . & . & . & . & . & . \\
0 & 0 & 0 & . & E_{N-2} & 1 & 0 \\
0 & 0 & 0 & . & 1 & E_{N-1} & 1 \\
0 & 0 & 0 & . & 0 & 1 & E_{N}
\end{array}\right\|
$$

contains level energies $E_{n}$ distributed according to (7.130). An explicit expression for the multidimensional Gaussian integral

$$
\begin{equation*}
\prod_{n=1}^{N}\left\{\int_{-\infty}^{\infty} \frac{d x_{n}}{\sqrt{\pi}}\right\} \exp \left[-i \sum_{n, m=1}^{N} x_{n} a_{n m} x_{m}\right]=\frac{1}{\sqrt{\operatorname{Det}\left|a_{n m}\right|}} \tag{7.133}
\end{equation*}
$$

valid for any matrix $a_{n m}$ with eigenvalues in the lower part of the complex plane, suggests a practical way to evaluate the ensemble averaged population (7.131). Indeed, taking into account the signs of the imaginary parts of $\varepsilon$ and $\xi$ for $a_{n m}=H_{n m}-\delta_{n, m} \varepsilon$ (or $\xi$ ), one finds

$$
\begin{align*}
& \rho_{N}(t)=\int_{-\infty-i \nu}^{\infty-i \nu} \frac{d \xi}{2 \pi} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \varepsilon}{2 \pi} \prod_{n=1}^{N}\left\{\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \frac{d x_{n} d \bar{x}_{n} d y_{n} d \bar{y}_{n}}{\pi^{2}}\right\} \\
& \exp \left[-i t(\varepsilon-\xi)-2 i \sum_{n=1}^{N-1}\left(x_{n} x_{n+1}+\bar{x}_{n} \bar{x}_{n+1}-y_{n} y_{n+1}-\bar{y}_{n} \bar{y}_{n+1}\right)\right. \\
& \left.-i \sum_{n=1}^{N}\left(\left(E_{n}-\varepsilon\right)\left(x_{n}^{2}+\bar{x}_{n}^{2}\right)-\left(E_{n}-\xi\right)\left(y_{n}^{2}+\bar{y}_{n}^{2}\right)\right)\right] \tag{7.134}
\end{align*}
$$

which after the ensemble average with (7.130) yields

$$
\begin{align*}
\rho_{N}(t)= & \int_{-\infty-i \nu}^{\infty-i \nu} \frac{d \xi}{2 \pi} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \varepsilon}{2 \pi} \prod_{n=1}^{N}\left\{\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \frac{d x_{n} d \bar{x}_{n} d y_{n} d \bar{y}_{n}}{\pi^{2}}\right\} \\
& \times \exp \left[-i t(\varepsilon-\xi)-2 i \sum_{n=1}^{N-1}\left(x_{n} x_{n+1}+y_{n} y_{n+1}-\bar{x}_{n} \bar{x}_{n+1}-\bar{y}_{n} \bar{y}_{n+1}\right)\right. \\
& \left.+\sum_{n=1}^{N}\left(i \varepsilon\left(x_{n}^{2}+\bar{x}_{n}^{2}\right)-i \xi\left(y_{n}^{2}+\bar{y}_{n}^{2}\right)-\frac{\gamma^{2}}{4}\left(x_{n}^{2}+\bar{x}_{n}^{2}-y_{n}^{2}-\bar{y}_{n}^{2}\right)^{2}\right)\right] . \tag{7.135}
\end{align*}
$$

These integrals are known in field theory as functional integrals on a grid and can be viewed as integrals over all trajectories $\vec{R}(n)$, given by all of the possible dependencies of the four-dimensional radius vector $\vec{R}=(x, \bar{x}, y, \bar{y})$ on the fictitious "time" $n$ that take discrete, integer, values. Note that one can also consider this four-dimensional real space as two-dimensional complex space of the variables $z=x+i \bar{x}$ and $z^{\prime}=y+i \bar{y}$ for which instead of (7.133) one has

$$
\begin{equation*}
\prod_{n=1}^{N}\left\{\int_{-\infty}^{\infty} \frac{d z_{n}^{*} d z_{n}}{\pi}\right\} \exp \left[-i \sum_{n, m=1}^{N} z_{n} a_{n m} z_{m}^{*}\right]=\frac{1}{\operatorname{Det}\left|a_{n m}\right|} \tag{7.136}
\end{equation*}
$$

which implies integration over all complex planes of each variable. One can also consider the complex plane $z$ as the phase space of a classical oscillator, where the square of the absolute value $\left|z_{n}\right|^{2}$ and the phase $\operatorname{Im} \ln z_{n}$ play the roles of action and phase, respectively.

## An Analogy between the Feynman Path Integral for a Boson and an Inverse Determinant

Prior to a detailed consideration of the integral (7.135), let us dwell on a simpler integral (7.133) for the case of $a_{n m}=H_{n m}-\delta_{n, m} \varepsilon$ and $H_{n m}$ of (7.132) with $E_{n} \equiv 0$, that is on the many-fold integral

$$
\begin{align*}
I_{N}(\varepsilon) & =\prod_{n=1}^{N}\left\{\int_{-\infty}^{\infty} \frac{d x_{n}}{\sqrt{\pi}}\right\} \exp \left[-2 i \sum_{n=1}^{N-1} x_{n} x_{n+1}+i \sum_{n=1}^{N} \varepsilon x_{n}^{2}\right] \\
& =\prod_{n=1}^{N}\left\{\int_{-\infty}^{\infty} \frac{d x_{n}}{\sqrt{\pi}}\right\} \exp \left[i \sum_{n=0}^{N}\left(x_{n}-x_{n+1}\right)^{2}+(\varepsilon-2) x_{n}^{2}\right], \tag{7.137}
\end{align*}
$$

where in the last line we have assumed $x_{0}=x_{N+1}=0$.
We trace an analogy between this expression and the Feynman path integral for a quantum particle. Indeed, considering $n$ as a continuous "time" variable and replacing finite differences by the first derivatives and the sums by the integrals, we obtain

$$
\begin{equation*}
I_{N} \simeq \int \exp \left\{i \int_{0}^{N+1}\left[\left(\frac{\partial x(n)}{\partial n}\right)^{2}+(\varepsilon-2) x^{2}(n)\right] d n\right\} \mathcal{D} x(n) \tag{7.138}
\end{equation*}
$$

with

$$
\begin{equation*}
x(N+1)=x(0)=0 . \tag{7.139}
\end{equation*}
$$

This is a standard path integral representation for the probability amplitude of a quantum particle with a classical Lagrangian

$$
\begin{equation*}
L\left(x, \frac{\partial x}{\partial n}\right)=\left(\frac{\partial x}{\partial n}\right)^{2}+(\varepsilon-2) x^{2} \tag{7.140}
\end{equation*}
$$

to remain in the initial position $x=0$ after a "time" $N+1$. This Lagrangian corresponds to a particle of mass $m=2$ moving in a parabolic potential $(2-\varepsilon) x^{2}$, whereas (7.138) simply represents the integral of $\exp (-i S)$ for the action $S=\int L\left(x, \frac{\partial x}{\partial n}\right) d n$ over all the trajectories starting at $x=0$ at $n=0$ and ending at the same point at $n=N+1$. The corresponding Hamiltonian

$$
\begin{equation*}
\widehat{H}=-\frac{1}{4}\left(\frac{\partial}{\partial x}\right)^{2}+(2-\varepsilon) x^{2} \tag{7.141}
\end{equation*}
$$

coincides with the Hamiltonian of a one-dimensional quantum harmonic oscillator of a "frequency" $\sqrt{4-2 \varepsilon}$. Note that the "time"-dependent Hamiltonian

$$
\begin{equation*}
\widehat{H}=\frac{1}{4} \widehat{p}^{2}+(2-\varepsilon) x^{2} \sum_{k=0}^{\infty} \delta(n-k) \tag{7.142}
\end{equation*}
$$

corresponds to the original integral (7.137) with the finite difference in the action sum. It employs the trajectories on a one-dimensional grid of integer


Fig. 7.9. The transmission probability amplitude for a set of $N$ levels, coupled by relay-like interaction, can be given in terms of Feynman trajectories on a discrete grid (a) and the corresponding potential curve (b) for a particle of mass $m=2$ which has probability amplitude to be in the initial state $x=0$ at the fictitious "time" $n=N$. Calculation of the population also requires integration over classical trajectories of spin $1 / 2$ (c).
numbers, as shown in Fig. 7.9. The parabolic dependence of the Hamiltonian (7.142) on the coordinate and momentum, suggests the assumption that the dynamics of the system is closely related to the dynamics of a harmonic oscillator.

## Bosonic Field, a Quantum Oscillator and Mapping over a Period

Earlier we have mentioned that the possibility to find an exact solution of (7.22) for a harmonic oscillator subjected to an arbitrary Raman field, results
from an underlying algebraic symmetry, summarized in (7.23). The same symmetry also exists for a harmonic oscillator with a time-dependent frequency, where the Hamiltonian (7.142) can be rewritten in terms of bosonic creation and annihilation operators $\widehat{a}^{\dagger}=c \widehat{x}+i f \widehat{p}$ and $\widehat{a}=c^{*} \widehat{x}-i f^{*} \widehat{p}$ :

$$
\begin{equation*}
\widehat{H}=-\frac{1}{16 f^{2}}\left(\widehat{a}-\widehat{a}^{\dagger}\right)^{2}+\frac{(2-\varepsilon)}{4 c^{2}}\left(\widehat{a}+\widehat{a}^{\dagger}\right)^{2} \sum_{k=0}^{\infty} \delta(n-k), \tag{7.143}
\end{equation*}
$$

in a form similar to the Hamiltonian in (7.21). A particular choice of the constants $c$ and $f$, constrained by the condition $\left[\widehat{a}, \widehat{a}^{\dagger}\right]=1$, is a matter of convenience. The symmetry manifests itself in the independence of the oscillation period on the oscillator amplitude, and is responsible for another important feature: the action $S$ corresponding to the Lagrangian (7.140) identically equals zero for all classical trajectories.

Though from the viewpoint of calculation, the path integral representation (7.137) with the allowance of (7.139) adds nothing new to the exact expression (7.133)

$$
\begin{align*}
I_{N}(\varepsilon) & =\prod_{n=1}^{N}\left\{\int_{-\infty}^{\infty} \frac{d x_{n}}{\sqrt{\pi}}\right\} \exp \left\{i \sum_{n=0}^{N}\left[\left(x_{n}-x_{n+1}\right)^{2}+(\varepsilon-2) x_{n}^{2}\right]\right\} \\
& =\operatorname{Det}^{-1 / 2}\left\|\begin{array}{cccc}
-\varepsilon & 1 & \cdot & 0 \\
1 & -\varepsilon & 0 \\
\cdot & \cdot & \cdot & \cdot \\
0 & 0 & \cdot & -\varepsilon
\end{array}\right\|=\frac{1}{\sqrt{\cos [N \arccos (-\varepsilon)]}}, \tag{7.144}
\end{align*}
$$

the analogy between the relay-like multilevel system and a free particle subjected to periodic $\delta$-like parabolic "kicks", suggests an additional intuitive insight, useful for the choice of an approximate description of disordered chains.

Note in this context that since integration with the kernel $e^{i\left(x_{n}-x_{n+1}\right)^{2}}$ is equivalent to the transformation $\exp \left\{-i \frac{1}{4} \widehat{p}^{2}\right\}$, and hence the operator

$$
\begin{equation*}
\widehat{U}=\exp \left\{-i(2-\varepsilon) x^{2}\right\} \exp \left\{-i \frac{1}{4} \widehat{p}^{2}\right\} \tag{7.145}
\end{equation*}
$$

gives the quantum evolution over a period for the particle, corresponding to the Hamiltonian (7.143). For the Heisenberg coordinate and momentum operators, the transformation (7.145) yields a mapping over a period

$$
\begin{align*}
& \widehat{U} \widehat{p} \widehat{U}^{-1}=\widehat{p}-2(2-\varepsilon) \widehat{x} \\
& \widehat{U} \widehat{x} \widehat{U}^{-1}=(-1+\varepsilon) \widehat{x}+\frac{1}{2} \widehat{p}, \tag{7.146}
\end{align*}
$$

which can also be written in the matrix form

$$
\binom{\widehat{p}}{\widehat{x}}_{n+1}=\left(\begin{array}{cc}
1 & -2(2-\varepsilon)  \tag{7.147}\\
\frac{1}{2} & -1+\varepsilon
\end{array}\right)\binom{\widehat{p}}{\widehat{x}}_{n} .
$$

The transposed matrix gives the change of the coefficients $r$ and $s$

$$
\binom{r}{s}_{n+1}=\left(\begin{array}{cc}
1 & \frac{1}{2}  \tag{7.148}\\
-2(2-\varepsilon) & -1+\varepsilon
\end{array}\right)\binom{r}{s}_{n} .
$$

for a linear combination $r \widehat{p}+s \widehat{x}$, as a result of the mapping. This $2 \times 2$ matrix has two eigenvalues $\lambda_{1,2}=\frac{\varepsilon}{2} \pm \sqrt{\frac{\varepsilon^{2}}{4}-1}$, such that $\left|\lambda_{1,2}\right|=1$ for $|\varepsilon|<2$, and the determinant of the matrix is equal to unity. For $\varepsilon=2 \sin \phi ; \pi / 2<\phi<$ $3 \pi / 2$ one has $\lambda_{1,2}=\exp [\mp i(\phi-\pi / 2)]$. The corresponding eigenvectors

$$
\begin{align*}
\widehat{a} & =\frac{1}{\sqrt{-\cos \phi}}\left[\widehat{p} \frac{e^{i(\phi+\pi / 2) / 2}}{2}+2 \widehat{x} \sin \left(\frac{\phi}{2}-\frac{\pi}{4}\right)\right] \\
\widehat{a}^{\dagger} & =\frac{1}{\sqrt{-\cos \phi}}\left[\widehat{p} \frac{e^{-i(\phi+\pi / 2) / 2}}{2}+2 \widehat{x} \sin \left(\frac{\phi}{2}-\frac{\pi}{4}\right)\right] \tag{7.149}
\end{align*}
$$

of the mapping (7.147) suggest a convenient choice of the bosonic creation and annihilation operators encountered in (7.143).

Mapping of $\widehat{a}$ and $\widehat{a}^{\dagger}$ is given by a simple multiplication $\widehat{a}=-i \exp (i \phi) \widehat{a}$, $\widehat{a}^{\dagger}=i \exp (-i \phi) \widehat{a}^{\dagger}$, whereas the state $|0\rangle$ corresponds for a positive $\phi$ to the wavefunction

$$
\begin{equation*}
\langle x \mid 0\rangle=\left|\frac{2 \cos \phi}{\pi}\right|^{1 / 4} e^{i x^{2}(1-i \exp (-i \phi))} \tag{7.150}
\end{equation*}
$$

which satisfies the relation $\widehat{a}|0\rangle=0$, thus being the ground state of the effective Hamiltonian $i \ln \widehat{U}$ of (7.145). In terms of the operators $\widehat{a}$ and $\widehat{a}^{\dagger}$ the operator $\widehat{U}$ of evolution over a period (7.145) adopts the form

$$
\begin{equation*}
\widehat{U}=\exp \left\{-i\left(\phi-\frac{\pi}{2}\right)\left(\widehat{a} \dagger \widehat{a}+\frac{1}{2}\right)\right\} \tag{7.151}
\end{equation*}
$$

For $\varepsilon=2 \sin \phi ; \pi / 2<\phi<3 \pi / 2$ one obtains the complex conjugate amplitudes (7.150) and inverse evolution. We note an explicit expression

$$
\begin{align*}
\widehat{x} & =\frac{1}{2 \sqrt{-\cos \phi}}\left(e^{i \tau} \widehat{a}^{\dagger}+e^{-i \tau} \widehat{a}\right) \\
\tau & =\phi / 2+\pi / 4 \tag{7.152}
\end{align*}
$$

for the coordinate operator in terms of the creation and the annihilation operators (7.149), where the coefficients in front of $\widehat{a}^{\dagger}$ and $\widehat{a}$ are just mutually conjugated phase factors.

## The Tunneling Transparency as Propagation of Four Bosons

The ensemble averaged transparency (7.135) can also be considered as a path integral for a "particle", by analogy to (7.138), although this particle now moves in a four-dimensional space and discrete time, having negative values
of mass for motion in the directions $\bar{x}$ and $\bar{y}$, whereas the problem itself does not have "oscillator" symmetry. The corresponding Hamiltonian

$$
\begin{align*}
\widehat{H}= & \frac{\widehat{p}_{x}^{2}}{4}+\frac{\widehat{p}_{\bar{x}}^{2}}{4}-\frac{\widehat{p}_{y}^{2}}{4}-\frac{\widehat{p}_{\bar{y}}^{2}}{4}+\frac{1}{4}\left[i \gamma^{2}\left(x^{2}+\bar{x}^{2}-y^{2}-\bar{y}^{2}\right)^{2}\right. \\
& \left.-4(2-\xi)\left(y^{2}+\bar{y}^{2}\right)+4(2-\varepsilon)\left(x^{2}+\bar{x}^{2}\right)\right] \sum_{k=0}^{\infty} \delta(n-k) \tag{7.153}
\end{align*}
$$

no longer has the same symmetry as (7.142), containing a bi-quadratic imaginary part of the "particle's" potential energy. This imaginary part results in an irreversible "decay", which implies that the rate of transition from the first to the last ( $N$-th) level of the chain decreases exponentially with the "time" $N$. The probability amplitude for the "particle" to be at the origin $x=y=\bar{x}=\bar{y}=0$ at "time" $N$, in the presence of the decay, gives, after inverse Fourier transformation, the population of the last level of the disordered chain.

Since at $\gamma=0,(7.153)$ is a sum of four independent Hamiltonians (7.142), for $\gamma \neq 0$, we can consider the problem in the interaction representation by just replacing coordinates in the interaction term $i \gamma^{2}\left(x^{2}+\bar{x}^{2}-y^{2}-\bar{y}^{2}\right)^{2} / 4$ by the unperturbed Heisenberg operators of these coordinates, that is

$$
\begin{equation*}
\widehat{H}_{i n t}=\sum_{k=0}^{\infty} \frac{i \gamma^{2}}{4}\left(x_{H}^{2}+\bar{x}_{H}^{2}-y_{H}^{2}-\bar{y}_{H}^{2}\right)^{2} \delta(n-k) \tag{7.154}
\end{equation*}
$$

According to (7.152) and the mapping relations on p. 349 for the creation and annihilation operators, the Heisenberg operators of coordinates have the form

$$
\begin{align*}
& \widehat{x}_{H}=\frac{1}{2 \sqrt{-\cos \phi}}\left(e^{i \tau}\left(-i e^{i \phi}\right)^{n} \widehat{a}^{\dagger}+e^{-i \tau}\left(i e^{-i \phi}\right)^{n} \widehat{a}\right), \\
& \widehat{\bar{x}}_{H}=\frac{1}{2 \sqrt{-\cos \phi}}\left(e^{i \tau}\left(-i e^{i \phi}\right)^{n} \widehat{b}^{\dagger}+e^{-i \tau}\left(i e^{-i \phi}\right)^{n} \widehat{b}\right), \\
& \widehat{y}_{H}=\frac{1}{2 \sqrt{-\cos \varphi}}\left(e^{i \tau^{\prime}}\left(i e^{-i \varphi}\right)^{n} \widehat{\bar{a}}^{\dagger}+e^{-i \tau^{\prime}}\left(-i e^{i \varphi}\right)^{n} \widehat{\bar{a}}\right), \\
& \widehat{\bar{y}}_{H}=\frac{1}{2 \sqrt{-\cos \varphi}}\left(e^{i \tau^{\prime}}\left(i e^{-i \varphi}\right)^{n} \widehat{\bar{b}}^{\dagger}+e^{-i \tau^{\prime}}\left(-i e^{i \varphi}\right)^{n} \widehat{\bar{b}}\right), \tag{7.155}
\end{align*}
$$

with $\xi=2 \sin \varphi ;-\pi / 2<\varphi<\pi / 2$, which results from the mapping over $n$ periods of all of the bosonic operators $\widehat{a}, \widehat{a}^{\dagger}, \widehat{b}, \widehat{b}^{\dagger}, \widehat{\bar{a}}, \widehat{\bar{a}}^{\dagger}, \widehat{\bar{b}}$, and $\widehat{\bar{b}}^{\dagger}$ corresponding to all four dimensions. The phase factor $e^{i \tau^{\prime}}$, implicitly given by (7.149), coincides with $e^{i \tau}$ after replacing $\phi$ by $\varphi$ and complex conjugation.

We concentrate on the case of $\gamma \ll 1$, since the opposite extreme corresponds to the case when the neighboring levels in the cascade are strongly off resonance and therefore their population is small, according to (7.89). In this limit we can consider an average of the Hamiltonian (7.154) over a large number of periods

$$
\begin{equation*}
\widehat{H}_{i n t}=\frac{i \gamma^{2}}{4}\left\langle\left(x_{H}^{2}+\bar{x}_{H}^{2}-y_{H}^{2}-\bar{y}_{H}^{2}\right)^{2}\right\rangle_{n} \tag{7.156}
\end{equation*}
$$

drop out all the terms oscillating with "time", that is the terms dependent on $\phi n$ or $\varphi n$, and concentrate only on the smallest eigenvalue of the remaining Hamiltonian, which corresponds to the vacuum state $|0\rangle$ of all four oscillators. One determines this value by making use of the explicit forms

$$
\begin{align*}
\widehat{x}_{H}^{2} & =\frac{1}{-4 \cos \phi}\left(e^{2 i \tau}\left(-i e^{i \phi}\right)^{2 n} \widehat{a}^{\dagger 2}+e^{-2 i \tau}\left(i e^{-i \phi}\right)^{2 n} \widehat{a}^{2}+\left\{\widehat{a}^{\dagger} \widehat{a}\right\}\right), \\
\left\langle\widehat{x}_{H}^{2}\right\rangle_{n} & =\frac{1}{-4 \cos \phi}\left\{\widehat{a}^{\dagger} \widehat{a}\right\} \xrightarrow{\langle 0| . .|0\rangle} \frac{1}{-4 \cos \phi} \\
\left\langle\widehat{\widehat{x}}_{H}^{4}\right\rangle_{n} & =\frac{1}{16 \cos ^{2} \phi}\left(\left\{\widehat{a}^{\dagger 2} \widehat{a}^{2}\right\}+\left\{\widehat{a}^{\dagger} \widehat{a}\right\}^{2}\right) \xrightarrow{\langle 0| . .|0\rangle} \frac{3}{16 \cos ^{2} \phi}, \tag{7.157}
\end{align*}
$$

of the Heisenberg coordinate operators and their averages in the ground state. Substituting into (7.156) the explicit mean values (7.157) and similar expressions for the other averaged operators, we find the smallest eigenvalue $i \gamma^{2}\left(\cos ^{-2} \phi+\cos ^{-2} \varphi-\cos ^{-1} \varphi \cos ^{-1} \phi\right) / 8 \sim i \gamma^{2} / 8 \cos \phi \cos \varphi$ which gives the slowest rate of "decay" with the course of "time" $n$. It therefore determines the asymptotic behavior at the limit of long chain $N \gg 1$.

The ground state dominates in the expression for the evolution operator in the interaction representation

$$
\begin{equation*}
\widehat{U}_{i n t}(n) \simeq|0\rangle e^{-\gamma^{2} n / 8 \cos \phi \cos \varphi}\langle 0|, \tag{7.158}
\end{equation*}
$$

and hence for the original problem we find an expression

$$
\begin{align*}
\rho_{N}(t)= & \int_{-\infty-i \nu}^{\infty-i \nu} \frac{d \xi}{2 \pi} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \varepsilon}{2 \pi} e^{-i t(\varepsilon-\xi)-\gamma^{2} N / 8 \cos \phi \cos \varphi}\left|\frac{2 \cos \varphi}{\pi} \frac{2 \cos \phi}{\pi}\right|^{1 / 2} \\
& \operatorname{Det}^{-1}\left\|\begin{array}{ccc}
-\varepsilon & 1 & \cdot \\
1 & -\varepsilon & 0 \\
\cdot & \cdot & \cdot \\
0 & 0 & \cdot \\
\hline
\end{array}\right\| \varepsilon^{\prime}\left\|\operatorname{Det}^{-1}\right\| \begin{array}{ccc}
-\xi & 1 & \cdot \\
1 & -\xi \cdot & 0 \\
. & \cdot & \cdot \\
0 & 0 & \cdot \\
\hline
\end{array} \| \tag{7.159}
\end{align*}
$$

where one recognizes the exponent corresponding to the evolution operator (7.158) in the interaction representation, the factors $\frac{2 \cos \varphi}{\pi}$ and $\frac{2 \cos \phi}{\pi}$ that according to (7.150) allow for the projection of the states $\langle 0|$ to the point $x_{0}=$ $y_{0}=\bar{x}_{0}=\bar{y}_{0}=0$ which is the origin of the trajectories, and the determinants that describe the free evolution in agreement with (7.144). The last diagonal terms $\varepsilon^{\prime}=1+\varepsilon-i \cos \phi$ and $\xi^{\prime}=1+\xi+i \cos \varphi$ in these determinants come from the projection of the evolution operators for unperturbed motion to the slowest decaying state $|0\rangle$, instead of setting $x_{N+1}=y_{N+1}=\bar{x}_{N+1}=\bar{y}_{N+1}=$ 0 which modifies the integral (7.137). Note that $\varepsilon=2 \sin \phi ; \xi=2 \sin \varphi$ as earlier.

Equation (7.159) has a clear physical meaning: the population of the last level of a disordered chain of $N$ levels has almost the same time behavior
as that of the infinite system of resonant levels depicted in Fig. 7.1(a) for $V=1$, becoming different from zero once the front of the population distribution $\left|\psi_{n}(t)\right|^{2}=J_{n}^{2}(2 t)$ of (7.11) approaches the level $n=N$ at the time $t \sim N / 2$. The disorder just results in attenuation of this distribution by the exponentially small factor $\sim e^{-\gamma^{2} N / 8}$ by analogy to the case of (7.97).

For the case of small $\gamma$, let us consider the long-time limit $\varepsilon \rightarrow \xi$ of the tunneling probability of a specific domain of spectrum at a given energy $\eta=(\varepsilon+\xi) / 2$ corresponding to a certain value of the density of states $g(\eta)$. For the case of the Lorentzian distribution (7.95) of the level energies, the density $g(\eta)$ coincides with $g(\varepsilon)$, given by (7.103) and shown in Fig. 7.6(d). The mean density of states for the Gaussian distribution (7.130) is clearly identical to (7.103) for small $\gamma$, and can be written in the form $g(\eta)=1 / \cos (\arcsin (\eta / 2))$. Therefore, in terms of the state density, the smallest eigenvalue of $\widehat{H}_{\text {int }}$ is $i \gamma^{2} g^{2}(\eta) / 8$, which takes the maximum value at the edges of the spectrum $\eta \rightarrow \pm 2$ where $g(\eta) \rightarrow \infty$. The maximum transparency comes from the center of the spectrum $\eta \rightarrow 0$, which is the spectrum domain responsible for the overall transparency $\sim e^{-\gamma^{2} N / 8}$. This is consistent with the dependence of the localization length on the energy depicted in Fig. 7.8(a) which clearly shows that the localization length assumes the minimum values at the spectrum edges.

### 7.5.2 An Analogy With the Liouville Equation

The different signs of the $x$ and $y$ variables in (7.153), and the imaginary interaction (7.154), suggest the employment of the density matrix formalism and the Liouville equation instead of the Hamiltonian formalism and the Schrödinger equation, as a more suitable technique. Indeed, let us consider the combination

$$
\begin{align*}
\widetilde{\rho}_{n}\left(x_{n}, \bar{x}_{n} ; y_{n}, \bar{y}_{n}\right)= & \prod_{k=1}^{n-1}\left\{\int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} \frac{d x_{k} d \bar{x}_{k} d y_{k} d \bar{y}_{k}}{\pi^{2}}\right\} \\
& \exp \left[-2 i \sum_{k=1}^{n-1}\left(x_{k} x_{k+1}+\bar{x}_{k} \bar{x}_{k+1}-y_{k} y_{k+1}-\bar{y}_{k} \bar{y}_{k+1}\right)\right. \\
& \sum_{k=1}^{n}\left(i \eta\left(x_{k}^{2}+\bar{x}_{k}^{2}-y_{k}^{2}-\bar{y}_{k}^{2}\right)+i \zeta\left(x_{k}^{2}+\bar{x}_{k}^{2}+y_{k}^{2}+\bar{y}_{k}^{2}\right)\right) \\
& \left.-\frac{\gamma^{2}}{4} \sum_{k=1}^{n}\left(x_{k}^{2}+\bar{x}_{k}^{2}-y_{k}^{2}-\bar{y}_{k}^{2}\right)^{2}\right] \tag{7.160}
\end{align*}
$$

entering (7.135), where we take the part of the integrand dependent only on $n$ first sets of variables and perform integrations over all variables except the last four: $x_{n}, \bar{x}_{n} ; y_{n}$, and $\bar{y}_{n}$. We identify this combination with the "density matrix", evolving in the course of fictitious "time" $n$ which we mark by the
tilde sign in order to distinguish between this fictitious object and the real density matrix. For this object one finds the recurrence relation

$$
\begin{align*}
& \tilde{\rho}_{n+1}(x, \bar{x} ; y, \bar{y})=\widehat{U}(x) \widehat{U}(\bar{x}) \widetilde{\rho}_{n}(x, \bar{x} ; y, \bar{y}) \widehat{U}^{*}(y) \widehat{U}^{*}(\bar{y}) \\
& \exp \left[i \zeta\left(x^{2}+\bar{x}^{2}+y^{2}+\bar{y}^{2}\right)-\frac{\gamma^{2}}{4}\left(x^{2}+\bar{x}^{2}-y^{2}-\bar{y}^{2}\right)^{2}\right] \tag{7.161}
\end{align*}
$$

with

$$
\begin{equation*}
\widehat{U}(x)=\exp \left\{-i(2-\eta) x^{2}\right\} \exp \left\{\frac{-i}{4} \widehat{p}_{x}^{2}\right\} \tag{7.162}
\end{equation*}
$$

by analogy to (7.145).
In the interaction representation

$$
\widetilde{\rho}_{n}(x, \bar{x} ; y, \bar{y})=(\widehat{U}(x) \widehat{U}(\bar{x}))^{n} \widetilde{\rho}_{n}^{(i n t)}(x, \bar{x} ; y, \bar{y})\left(\widehat{U}^{*}(y) \widehat{U}^{*}(\bar{y})\right)^{n}
$$

and for small $\gamma$ and $\zeta$ one finds by a Taylor expansion in $\gamma^{2}$

$$
\begin{align*}
& \widetilde{\rho}_{n+1}^{(\text {int })}(x, \bar{x} ; y, \bar{y})-\widetilde{\rho}_{n}^{(\text {int })}(x, \bar{x} ; y, \bar{y})=-\frac{\gamma^{2}}{2} \widehat{Z}_{H}^{2} \widetilde{\rho}_{n}^{(\text {int })}(x, \bar{x} ; y, \bar{y}){\widehat{Z^{\prime}}}_{H}^{2} \\
& +\left(i \zeta \widehat{Z}_{H}^{2}-\frac{\gamma^{2}}{4} \widehat{Z}_{H}^{4}\right) \widetilde{\rho}_{n}^{(\text {int })}(x, \bar{x} ; y, \bar{y})+\widetilde{\rho}_{n}^{(\text {int })}(x, \bar{x} ; y, \bar{y})\left(i \zeta{\widehat{Z^{\prime}}}_{H}^{2}-\frac{\gamma^{2}}{4}{\widehat{Z^{\prime}}}_{H}^{4}\right) \tag{7.163}
\end{align*}
$$

where $\widehat{Z}_{H}^{2}=\left(\widehat{x}_{H}^{2}+\widehat{\bar{x}}_{H}^{2}\right), \widehat{Z}_{H}^{2}=\left(\widehat{y}_{H}^{2}+\widehat{\bar{y}}_{H}^{2}\right)$ are the left and right Heisenberg operators, respectively. Equations (7.155) and (7.157) specify these operators and their powers in terms of the bosonic creation and annihilation operators (7.149), where $\varepsilon$ or $\xi$ should be replaced by $\eta$. To be consistent within the approximation of small $\gamma$ and $\zeta$, we should also replace the finite difference with the "time" derivative $\partial / \partial n$ and take averages of the operators over many periods. This yields

$$
\begin{align*}
\frac{\partial \widetilde{\rho}_{n}^{(i n t)}}{\partial n}= & i \zeta\left\{\widehat{H}_{o}, \widetilde{\rho}_{n}^{(\text {int })}\right\}-\frac{\gamma^{2}}{4}\left[\widehat{H}_{o},\left[\widehat{H}_{o}, \widetilde{\rho}_{n}^{(\text {int })}\right]\right] \\
& -\frac{\gamma^{2}}{4}\left[\widehat{A}_{2}^{\dagger},\left[\widehat{A}_{2}, \widetilde{\rho}_{n}^{(\text {int })}\right]\right]-\frac{\gamma^{2}}{4}\left[\widehat{A}_{2},\left[\widehat{A}_{2}^{\dagger}, \widetilde{\rho}_{n}^{(i n t)}\right]\right] \tag{7.164}
\end{align*}
$$

where [, ] and $\{$,$\} denote commutator and anticommutator, respectively, and$

$$
\begin{align*}
\widehat{H}_{o} & =\omega\left(\widehat{a}^{\dagger} \widehat{a}+\widehat{b}^{\dagger} \widehat{b}+1\right) ; \widehat{A}_{2}^{\dagger}=\omega\left(\widehat{a}^{\dagger 2}+\widehat{b}^{\dagger 2}\right) \\
\widehat{A}_{2} & =\omega\left(\widehat{a}^{2}+\widehat{b}^{2}\right) ; \omega=1 / 4 \cos (\arcsin \eta / 2) \tag{7.165}
\end{align*}
$$

The operators $b$ and $b^{\dagger}$ correspond here to the coordinates with bars.
Equation (7.164) resembles the Liouville equation, where the relaxation term governing the evolution of the density matrix in the interaction representation originates from the ensemble averaged disorder $\gamma^{2}$. It is an analog of the relaxation operator for the regular density matrix of an open quantum system, originating from random time-dependent external interventions into
the coherent dynamics. Note that after substitution of the explicit expressions for the commutators and the anticommutator, the relaxation operator contains three different components $\widehat{R}_{l l}, \widehat{R}_{r r}$, and $\widehat{R}_{l r}$ that act on $\widetilde{\rho}_{n}^{(\text {int })}$ from the left-hand side, from the right-hand side, and from both sides, respectively. These relaxation operators are linear, although their action on the density matrix cannot always be written as a matrix multiplication as it was discussed in the context of (5.58). We therefore employ the notations with double "hats" $\widehat{\widehat{R}}_{l l}, \widehat{\widehat{R}}_{r r}$, and $\widehat{\widehat{R}}_{l r}$ on top of these operators in the situation where this circumstance is important.

For large $n$, the dominating contribution comes from the eigenvector (in this case of the Liouville equation it is the eigen "density" matrix) which corresponds to the eigenvalue of the Liouville operator with the smallest imaginary part. The "density" matrix satisfying (7.164) is a diagonal matrix in the oscillator quantum numbers representation $\sum_{m k}|m k\rangle \widetilde{\rho}_{n}^{(m k)}\langle k m|$. The main contribution to the slowest decaying term comes from the vacuum state $m=k=0$ for which the Liouville equation reads

$$
\begin{equation*}
\frac{\partial \widetilde{\rho}_{n}^{(0,0)}}{\partial n}=i 2 \zeta \omega \widetilde{\rho}_{n}^{(0,0)}-2(\gamma \omega)^{2} \widetilde{\rho}_{n}^{(0,0)} \tag{7.166}
\end{equation*}
$$

The dominating contribution therefore reads

$$
\widetilde{\rho}_{n}^{(\text {int })}(x, \bar{x} ; y, \bar{y})=\langle x, \bar{x} \mid 0\rangle \widetilde{\rho}_{n}^{(0,0)}\langle 0 \mid y, \bar{y}\rangle,
$$

where the coordinate representation of the wavefunctions is given explicitly in (7.150). We have thus reproduced the result of the previous subsection (p. 352) and have obtained for $\zeta \rightarrow 0$ the static transparency decreasing as $\sim e^{-\gamma^{2} N / 8}$ in the center of the band $(\eta=0)$ where the "frequency" assumes, according to (7.165), the minimum value $\omega=1 / 4$.

### 7.5.3 Classical Fermionic Fields for the Population Dynamics

We now consider the evolution of the population $\rho_{n}(t)$ in an infinite relay of levels of energies $E_{n}$, randomly distributed around the mean position $\Omega=$ 0 with Gaussian statistics (7.130). We calculate $\rho_{n}(t)$ for the case where all of the population of the chain has been initially localized at level 0 . In other words, the initial condition $\psi_{0}(t=0)=1$ holds for the Schrödinger equation (7.1), where the index $n$ assumes both positive and negative integer values. The constant coupling $V_{n}$ is taken as the energy unit. From (7.93) one immediately finds the population

$$
\begin{equation*}
\rho_{n}(t)=\int_{-\infty-i \nu}^{\infty-i \nu} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \xi d \varepsilon}{4 \pi^{2}} e^{-i(\varepsilon-\xi) t} \frac{M_{0}^{+}(\varepsilon) M_{n}^{-}(\varepsilon)}{\operatorname{Det}|\widehat{H}-\varepsilon|} \frac{M_{0}^{+}(\xi) M_{n}^{-}(\xi)}{\operatorname{Det}|\widehat{H}-\xi|} \tag{7.167}
\end{equation*}
$$

Though this expression resembles (7.131) for the tunneling transparency, it contains an important additional factor in the numerator of the integrand
given by the product of four minors of the matrices $(\widehat{H}-\varepsilon)_{n, n^{\prime}}$ and $(\widehat{H}-\xi)_{n, n^{\prime}}$. As earlier, one can represent the determinants in the denominator with the help of the functional integrals (7.133). We are now going to introduce a technique which allows one to represent the factors in the numerator also in the form of functional integrals, and thereby to perform the Gaussian ensemble average by analogy to (7.134) and (7.135).

## Analogy Between Minors and an Ensemble of Spins 1/2

As we have already seen in the previous section, the functional integral (7.137) for the inverse square root of the determinant can be interpreted in terms of a one-dimensional particle, with a Hamiltonian (7.142) moving in a parabolic potential, periodically changing in time. The integral equals the probability amplitude for the particle to remain at the origin after $N+1$ periods, whereas the evolution after one period is given by (7.145). We can build a similar construction for the minors starting from the recurrence relation (7.43). However, now the representing particle does not correspond to an infinite-dimensional Hilbert space of a harmonic oscillator, but to just twodimensional Hilbert space of a spin $1 / 2$. Indeed, by introducing "spinors $1 / 2$ " given as $\left\{M_{n}^{+}, i M_{n-1}^{+}\right\}$one can rewrite (7.43) for $V_{n} \equiv 1$ in the matrix form

$$
\left\{\begin{array}{c}
M_{n+1}^{+}  \tag{7.168}\\
i M_{n}^{+}
\end{array}\right\}=\left(\begin{array}{cc}
E_{n}-\varepsilon i \\
i & 0
\end{array}\right)\left\{\begin{array}{c}
M_{n}^{+} \\
i M_{n-1}^{+}
\end{array}\right\}
$$

where the $2 \times 2$ matrix describes the evolution over the period by analogy to (7.145).

The upper diagonal matrix element $(s=1 / 2)$ of the product of these matrices corresponding to the evolution over $N$ periods yields the entire determinant Det $|\widehat{H}-\varepsilon|$, whereas that for $n+1$ periods yields $M_{n}^{+}$. In other words, projection to the upper state of the "spin" plays the same role as the projection to the state $|x=0\rangle$ of the oscillator. Note that the operator of the evolution over a period

$$
\left(\begin{array}{rr}
E_{n}-\varepsilon & i  \tag{7.169}\\
i & 0
\end{array}\right)=\left(E_{n}-\varepsilon\right) \widehat{\sigma}^{+} \widehat{\sigma}+i \widehat{\sigma}^{+}+i \widehat{\sigma}
$$

can be given in terms of the spin raising and lowering operators $\widehat{\sigma}^{+}$and $\widehat{\sigma}$. One can introduce "trajectories" (see Fig. 7.9c) for the spin $1 / 2$ that evidently are discrete, not only in "time" $n$ but also in "space", since the spin assumes only two values $s_{ \pm}= \pm 1 / 2$. The trajectories are just sequences $\left\{\left(s_{ \pm}\right)_{n}\right\}$ of $s_{+}$or $s_{-}$for each $n$, and each trajectory corresponds unambiguously to one of $2^{N}$ terms in the explicit representation of the three-diagonal determinant Det $|\widehat{H}-\varepsilon|$. We assume the standard vector notation

$$
\begin{array}{ll}
\left|s_{+}\right\rangle=\binom{1}{0} ; & \left|s_{-}\right\rangle=\binom{0}{1} \\
\left\langle s_{+}\right|=\left(\begin{array}{ll}
1 & 0
\end{array}\right) ; & \left\langle s_{-}\right|=\left(\begin{array}{ll}
0 & 1
\end{array}\right) \tag{7.170}
\end{array}
$$

for the spin eigenstates.
The correspondence between the terms of the determinant and the spin trajectoties can be established as follows. We develop the determinant over the elements of the last column. The diagonal matrix element then corresponds to $s_{+}$and has the factor $\left(E_{1}-\varepsilon\right)$ in front of the minor $M_{1}^{+}$, whereas the off-diagonal matrix element corresponds to $s_{-}$and gives the unity factor in front of the minor $\left(M_{2}^{+}\right)_{a p}$. The latter results from appending the minor $M_{2}^{+}$by one line and one column, containing zeros everywhere apart from unity on the diagonal position. The same development applied to $M_{1}^{+}$ or $\left(M_{2}^{+}\right)_{a p}$ provides us with the second element of the trajectory, and the procedure can be repeated up to the end, thus yielding the representation of the determinant as a sum of different products of unities and the factors $\left(E_{n}-\varepsilon\right)$, whereas each of the products corresponds to a different trajectory.

We illustrate this for the simplest example of a $2 \times 2$ determinant. On the one hand, straightforward calculation via expansion over the elements of the last column yields

$$
\begin{align*}
\operatorname{Det}\left|\begin{array}{cc}
E_{2}-\varepsilon & 1 \\
1 & E_{1}-\varepsilon
\end{array}\right| & =M_{1}^{+}\left(E_{1}-\varepsilon\right)-\left(M_{2}^{+}\right)_{a p} 1 \\
& =\left(E_{2}-\varepsilon\right)\left(E_{1}-\varepsilon\right)-1, \tag{7.171}
\end{align*}
$$

where the $1 \times 1$ minors are $M_{1}^{+}=\left(E_{2}-\varepsilon\right)$ and $\left(M_{2}^{+}\right)_{a p}=(1)$. On the other hand, with the help of (7.168) one obtains

$$
\begin{align*}
\operatorname{Det}|\widehat{H}-\varepsilon|= & \left\langle s_{+}\right|\left(\begin{array}{rr}
E_{2}-\varepsilon & i \\
i & 0
\end{array}\right)\left(\begin{array}{rr}
E_{1}-\varepsilon & i \\
i & 0
\end{array}\right)\left|s_{+}\right\rangle \\
= & \left\langle s_{+}\right|\left(\begin{array}{rr}
E_{2}-\varepsilon & i \\
i & 0
\end{array}\right)\left|s_{+}\right\rangle\left\langle s_{+}\right|\left(\begin{array}{rr}
E_{1}-\varepsilon & i \\
i & 0
\end{array}\right)\left|s_{+}\right\rangle \\
& +\left\langle s_{+}\right|\left(\begin{array}{rr}
E_{2}-\varepsilon & i \\
i & 0
\end{array}\right)\left|s_{-}\right\rangle\left\langle s_{-}\right|\left(\begin{array}{rr}
E_{1}-\varepsilon & i \\
i & 0
\end{array}\right)\left|s_{+}\right\rangle, \tag{7.172}
\end{align*}
$$

where in agreement with (7.169,7.170), the last two terms show the contribution $\left(E_{2}-\varepsilon\right)_{2}\left\langle s_{+}\right| \hat{\sigma}^{+} \widehat{\sigma}\left|s_{+}\right\rangle_{2}\left(E_{1}-\varepsilon\right)_{1}\left\langle s_{+}\right| \hat{\sigma}^{+} \widehat{\sigma}\left|s_{+}\right\rangle_{1}$ of the trajectory $\left(s_{+}\right)_{2}\left(s_{+}\right)_{1}$ and the contribution ${ }_{1}\left\langle s_{+}\right| \widehat{\sigma}^{+}\left|s_{-}\right\rangle_{2}(-1)_{2}\left\langle s_{-}\right| \widehat{\sigma}\left|s_{+}\right\rangle_{1}$ of the trajectory $\left(s_{-}\right)_{2}\left(s_{+}\right)_{1}$. The indices $n$ of $\left|s_{ \pm}\right\rangle_{n}$ enumerate columns of the determinant, whereas that of ${ }_{n}\left\langle s_{ \pm}\right|$enumerate the lines. Note that for a higher order determinant, the trajectory $\left(s_{r}\right)_{N}\left(s_{u}\right)_{c} \ldots\left(s_{s}\right)_{b}\left(s_{v}\right)_{a}\left(s_{w}\right)_{1}$ implies that the operators (7.169) are sequentially placed in between ${ }_{a}\left\langle s_{v}\right|$ and $\left|s_{w}\right\rangle_{1}$, in between ${ }_{b}\left\langle s_{s}\right|$ and $\left|s_{v}\right\rangle_{a}$, etc. until one finally reaches the last operator in between ${ }_{N}\left\langle s_{r}\right|$ and $\left|s_{u}\right\rangle_{c}$, at the end of the trajectory. The running indic7.177es $a=1,2 ; b=1,2 ; c=1,2$ mark the trajectory.

## Classical Anticommuting Variables

The main problem in describing a quantum system, with the help of a functional integral over classical trajectories, is finding a proper "action" functional $\exp [S]$ and the corresponding classical integration variables. For the particular problem under consideration, the functional integration should be equivalent to summation over the "spin" trajectories. Each "spin" trajectory $\left\{\left(s_{ \pm}\right)_{n}\right\}$ corresponds to a spin eigenstate of an ensemble of $N$ spins $1 / 2$, explicitly given by the tensor product $\prod_{n=1}^{N} \otimes\left|s_{ \pm}\right\rangle_{n}$ of $N$ individual spin states (either $s_{+}$or $s_{-}$) enumerated by the index $n$, such that the trajectory sum can also be considered as a trace over these eigenstates. However, the functional integrals convenient for the description of spin ensembles do not employ spins as the coordinates of classical trajectories, but rely on a different construction.

The fact that the problem is formulated in terms of the ensemble of spins of $1 / 2$ suggests the use of the anticommuting fermionic operators. In the context of tunneling transparency, we have seen the correspondence between a quantum harmonic oscillator described by the bosonic operators $\widehat{a}^{\dagger}$ and $\widehat{a}$ with the commutation relation $\left[\hat{a}, \widehat{a}^{\dagger}\right]=1$ and the functional integral over the regular (commuting) classical coordinates $x_{n}, \bar{x}_{n}$. By analogy, classical anticommuting coordinates $\chi_{n}$ and $\bar{\chi}_{n}$ with the anticommutators $\left\{\chi_{n}, \bar{\chi}_{m}\right\}=$ $0,\left\{\chi_{n}, \chi_{m}\right\}=0$, and $\left\{\bar{\chi}_{n}, \bar{\chi}_{m}\right\}=0$ are related to the quantum spin operators $\hat{\sigma}^{+}$and $\widehat{\sigma}$ of (7.169) satisfying the anticommutation relation $\left\{\widehat{\sigma}^{+}, \widehat{\sigma}\right\}=1$. These classical anticommuting coordinates are known as nilpotent, graded, or Grassmann variables. The anticommutation relations imply that $\chi_{n}^{2}=0$ and $\bar{\chi}_{n}^{2}=0$. Moreover, they also imply that each anticommuting coordinate enters, at most once, to any non-vanishing product of the coordinates, and in particular $\chi_{n} \bar{\chi}_{n} \chi_{n} \bar{\chi}_{n}=0$.

A certain physical meaning can be attributed to the variables $\chi_{n}$ and $\bar{\chi}_{n}$. The variable $\bar{\chi}_{n}$ can be interpreted as the creation of a "particle" at the position $n$, whereas $\chi_{n}$ denotes either annihilation of the existing "particle" or creation of an "antiparticle" at the same position. The combination $\bar{\chi}_{n} \chi_{n}$ may thus be interpreted as the number of particles ( 0 or 1 ) at the site $n$, whereas the combination $\bar{\chi}_{n+1} \chi_{n}$ indicates creation of a "particle-antiparticle" couple at the sites $n$ and ( $n+1$ ), respectively. Each "particle" or antiparticle can be created and annihilated only once, as follows from the anticommutation relations for the fermionic operators.

Since each "particle" has "spin" $1 / 2$, creation or annihilation of a couple of "particles" changes the total "spin" by unity. In the case of the nearest neighbor, relay-like interaction, resulting in three-term recurrence relations for the minors, and in accordance with (7.137) following from these relations, the total "spin" always remains $1 / 2$, whereas creation or annihilation of a couple just corresponds to "spin flip". We note that the "pair" corresponding to spin 1, and its creation and annihilation operators given by products of two Grassmann coordinates (bi-Grassmann), are bosonic operators commut-
ing among themselves. The operators of "pairs" also commute with single "particle" operators. Therefore one can manipulate with these operators as regular numbers.

Now each term of the determinant, which corresponds to a certain choice of the matrix elements in the sequential development of the minors over the elements of the last column, can also be attributed to a "spin" trajectory, and each of the "spin" trajectories corresponds in turn to a certain sequence of creation and annihilation of "particles" and "pairs". Let us illustrate this for the simplest example of (7.171), (7.172)

$$
\begin{align*}
\operatorname{Det}\left|\begin{array}{cc}
E_{2}-\varepsilon & 1 \\
1 & E_{1}-\varepsilon
\end{array}\right| & =\left(E_{2}-\varepsilon\right)\left(E_{1}-\varepsilon\right)-1 ; \\
\left(E_{2}-\varepsilon\right)\left(E_{1}-\varepsilon\right) & \rightarrow\left(s_{+}\right)_{2}\left(s_{+}\right)_{1} \rightarrow \bar{\chi}_{2} \chi_{2} \bar{\chi}_{1} \chi_{1} ; \\
-1 \rightarrow\left(s_{-}\right)_{2}\left(s_{+}\right)_{1} & \rightarrow \bar{\chi}_{1} \chi_{2} \bar{\chi}_{2} \chi_{1} . \tag{7.173}
\end{align*}
$$

The first term of the determinant resulting from the choice of the matrix element $\left(E_{1}-\varepsilon\right)$ in the last column corresponds to the trajectory $\left(s_{+}\right)_{2}\left(s_{+}\right)_{1}$ and to the creation and annihilation of "particle" 1 followed by the creation and annihilation of "particle" 2 , whereas the second term resulting from the choice of the matrix element 1 in the last column corresponds to the creation of a "pair" consisting of a "particle" at site 1 and an "antiparticle" at site 2 followed by their annihilation. Our next step is to find an algebraic way to reveal this correspondence.

## Functional Integrals over Anticommuting Variables

Consider a bilinear form

$$
\begin{align*}
& \overleftarrow{\bar{\chi}}(\widehat{H}-\varepsilon) \vec{\chi} \equiv\left(\bar{\chi}_{2} \bar{\chi}_{1}\right)\left(\begin{array}{cc}
E_{2}-\varepsilon & 1 \\
1 & E_{1}-\varepsilon
\end{array}\right)\binom{\chi_{2}}{\chi_{1}}  \tag{7.174}\\
= & \left(E_{2}-\varepsilon\right) \bar{\chi}_{2} \chi_{2}+\left(E_{1}-\varepsilon\right) \bar{\chi}_{1} \chi_{1}+\bar{\chi}_{2} \chi_{1}+\bar{\chi}_{1} \chi_{2} \equiv-S
\end{align*}
$$

and calculate the exponent $\exp [S]$ of this quantity, casting it in a Taylor series

$$
\begin{align*}
\exp [S]= & 1-\left(E_{2}-\varepsilon\right) \bar{\chi}_{2} \chi_{2}-\left(E_{1}-\varepsilon\right) \bar{\chi}_{1} \chi_{1}-\bar{\chi}_{2} \chi_{1}-\bar{\chi}_{1} \chi_{2} \\
& +\left(E_{2}-\varepsilon\right) \bar{\chi}_{2} \chi_{2}\left(E_{1}-\varepsilon\right) \bar{\chi}_{1} \chi_{1}+\bar{\chi}_{2} \chi_{1} \bar{\chi}_{1} \chi_{2} \tag{7.175}
\end{align*}
$$

where we retain only the terms linear in each coordinate, and allow for the fact that all binary products of the variables commute among themselves. Let us concentrate only on the two last terms of (7.175) containing all of the fermionic coordinates, and with the help of the anticommutation relations, make the order of the variables of the last term identical to those of the preceding term. This yields

$$
\begin{equation*}
\exp [S]=\ldots+\left(E_{2}-\varepsilon\right)\left(E_{1}-\varepsilon\right) \bar{\chi}_{2} \chi_{2} \bar{\chi}_{1} \chi_{1}-\bar{\chi}_{2} \chi_{2} \bar{\chi}_{1} \chi_{1}, \tag{7.176}
\end{equation*}
$$

where, in the numerical factor in front of the coordinate product $\bar{\chi}_{2} \chi_{2} \bar{\chi}_{1} \chi_{1}$, one recognizes the determinant (7.173).

The functional $S$ can serve as the action of the path integrals over anticommuting variables, provided this integration is defined as a procedure ensuring the selection of all of the terms in the Taylor expansion that only contain products of all of the coordinates. In other words, a "particle" should be created and annihilated at each site $n$. This goal is met by following the definitions of the integrals over the Grassmannian variables

$$
\begin{equation*}
\int d \chi_{n}=0, \quad \int \chi_{n} d \chi_{n}=1 \tag{7.177}
\end{equation*}
$$

where the differentials of the fermionic variables also anticommute. Equation (7.177) together with the anticommutation relations

$$
\begin{align*}
& \left\{d \chi_{n} d \chi_{m}\right\}=\left\{d \chi_{n} d \bar{\chi}_{m}\right\}=0 \\
& \left\{\chi_{n}, \bar{\chi}_{m}\right\}=0 ; \quad\left\{\chi_{n}, \chi_{m}\right\}=0 ; \quad\left\{\bar{\chi}_{n}, \bar{\chi}_{m}\right\}=0 \tag{7.178}
\end{align*}
$$

can serve as an axiomatic definition of the Grassmann variables.
We therefore have

$$
\begin{equation*}
\int \exp [S] d \chi_{1} d \bar{\chi}_{1} d \chi_{2} d \bar{\chi}_{2}=\left(E_{2}-\varepsilon\right)\left(E_{1}-\varepsilon\right)-1 \tag{7.179}
\end{equation*}
$$

whereas the functional integral

$$
\begin{equation*}
\int \exp \left[-\sum_{n, m=1}^{N} \bar{\chi}_{n} a_{n m} \chi_{m}\right] \prod_{n=1}^{N}\left\{d \chi_{n} d \bar{\chi}_{n}\right\}=\operatorname{Det}\left|a_{n m}\right| \tag{7.180}
\end{equation*}
$$

is the generalization of (7.179) for the case of an arbitrary order $N$ of the determinant, and is an analog of (7.136) for the regular variable. This expression is valid for the general matrix $a_{n m}$ even when the three-term recurrence relations (7.43) and (7.168) do not hold. Note that for a minor $M_{k l}$ resulting from the elimination of the $k$-th line and $l$-th column from the matrix $a_{n m}$ one has an integral representation

$$
\begin{equation*}
M_{k l}=\int \bar{\chi}_{k} \chi_{l} \exp \left[-\sum_{n, m=1}^{N} \bar{\chi}_{n} a_{n m} \chi_{m}\right] \prod_{n=1}^{N}\left\{d \chi_{n} d \bar{\chi}_{n}\right\} . \tag{7.181}
\end{equation*}
$$

Two coordinates in the pre-exponential factor eliminate all of the terms containing the matrix elements of the $k$-th line and $l$-th column from the Taylor expansion of the integrand.

### 7.6 Population Dynamics in a Disordered Chain

With the help of Grassmann notation, we now consider the population of an infinite relay of levels with disordered energies given by (7.167). We note that for the relay-like systems with $V_{n}=1$ one has $M_{0}^{+} M_{n}^{-}=M_{0 n}$ and hence

$$
\begin{align*}
\rho_{n}(t)= & \int_{-\infty-i \nu}^{\infty-i \nu} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \xi d \varepsilon}{4 \pi^{2}} e^{-i(\varepsilon-\xi) t} \prod_{k=-\infty}^{\infty}\left\{\int d \chi_{k} d \bar{\chi}_{k} d \chi_{k}^{\prime} d \bar{\chi}_{k}^{\prime}\right\} \\
& \bar{\chi}_{0}^{\prime} \chi_{n}^{\prime} \bar{\chi}_{0} \chi_{n} \exp \left[-\sum_{k, m=-\infty}^{\infty}\left(\bar{\chi}_{k} a_{k m} \chi_{m}-\bar{\chi}_{k}^{\prime} a_{k m}^{\prime} \chi_{m}^{\prime}\right)\right] \\
& \prod_{k=-\infty}^{\infty}\left\{\int_{-\infty}^{\infty} \frac{d z_{k}^{*} d z_{k} d z_{k}^{\prime *} d z_{k}^{\prime}}{4 \pi^{2}}\right\} \exp \left[-i \sum_{k, m=-\infty}^{\infty}\left(z_{k} a_{k m} z_{m}^{*}-z_{k}^{\prime} a_{k m}^{\prime} z_{m}^{\prime *}\right)\right] \tag{7.182}
\end{align*}
$$

where $a_{k m}=\left(E_{k}-\varepsilon\right) \delta_{k, m}+\delta_{k, m+1}+\delta_{k, m-1}$ and $a_{k m}^{\prime}=\left(E_{k}-\xi\right) \delta_{k, m}+$ $\delta_{k, m+1}+\delta_{k, m-1}$. The integrals over the variables $z^{*}, z, z^{\prime *}$, and $z^{\prime}$ are just another form of the integral (7.134), given in the notation of (7.136), whereas the integrals over the anticommuting variables of the type (7.181), take into account the minors in the numerators of (7.167).

By analogy to (7.135), for the Gaussian distribution function (7.130) centered around $\Omega=0$, one finds the ensemble average of (7.182)

$$
\begin{align*}
\rho_{n}(t)= & \int_{-\infty-i \nu}^{\infty-i \nu} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \xi d \varepsilon}{4 \pi^{2}} \prod_{k=-\infty}^{\infty}\left\{\int_{-\infty}^{\infty} \frac{d z_{k}^{*} d z_{k} d z_{k}^{\prime *} d z_{k}^{\prime}}{4 \pi^{2}}\right\} \\
& \prod_{k=-\infty}^{\infty}\left\{\int d \chi_{k} d \bar{\chi}_{k} d \chi_{k}^{\prime} d \bar{\chi}_{k}^{\prime}\right\} \bar{\chi}_{0}^{\prime} \chi_{n}^{\prime} \bar{\chi}_{0} \chi_{n} e^{-i(\varepsilon-\xi) t+S} ; \\
S= & \sum_{k=-\infty}^{\infty}\left[\varepsilon \bar{\chi}_{k} \chi_{k}-\xi \bar{\chi}_{k}^{\prime} \chi_{k}^{\prime}+i \varepsilon\left|z_{k}\right|^{2}-i \xi\left|z_{k}^{\prime}\right|^{2}+\bar{\chi}_{k}^{\prime} \chi_{k \pm 1}^{\prime}-\bar{\chi}_{k} \chi_{k \pm 1}\right. \\
& \left.+i z_{k}^{\prime} z_{k \pm 1}^{\prime *}-i z_{k} z_{k \pm 1}^{*}-\frac{\gamma^{2}}{4}\left(i \bar{\chi}_{k}^{\prime} \chi_{k}^{\prime}-i \bar{\chi}_{k} \chi_{k}+\left|z_{k}\right|^{2}-\left|z_{k}^{\prime}\right|^{2}\right)^{2}\right] \tag{7.183}
\end{align*}
$$

where the subscript $k \pm 1$ denotes the sum of the terms with indices $k-1$ and $k+1$.

### 7.6.1 Population Dynamics and Propagating Fictitious Particles

Considering the inverse square root of a three-diagonal determinant in the context of tunneling transparency as a function of its order $n$, we have seen
that for $\gamma=0$, the functional integral (7.137) representing this determinant can be interpreted as the probability amplitude of a periodically driven harmonic oscillator evolving in a fictitious "time" $n$. The evolution over a period (7.151) written with the help of the creation and the annihilation operators (7.149) suggests the employment of the oscillator energy eigenfunctions, where the most important role is played by the ground state (7.150). We have expressed the tunneling transparency in terms of four oscillators composing a "density matrix" (7.160) and considering $\gamma^{2} \neq 0$ as a perturbation. Now we employ the same idea for the anticommuting variables.

First we separate the "fermionic" part of (7.183) and introduce for $\gamma=0$ the "fermionic density matrix" $\bar{\rho}_{n}\left(\bar{\chi}_{k}, \chi_{k}, \bar{\chi}_{k}^{\prime}, \chi_{k}^{\prime}\right)$

$$
\begin{align*}
& \bar{\rho}_{n}\left(\bar{\chi}_{n}, \chi_{n}, \bar{\chi}_{n}^{\prime}, \chi_{n}^{\prime}\right)=\exp \left(\varepsilon \bar{\chi}_{n} \chi_{n}-\xi \bar{\chi}_{n}^{\prime} \chi_{n}^{\prime}\right) \prod_{k=-\infty}^{n-1}\left\{\int d \chi_{k} d \bar{\chi}_{k} d \chi_{k}^{\prime} d \bar{\chi}_{k}^{\prime}\right\} \\
& \exp \left[\sum_{k=-\infty}^{n-1} \varepsilon \bar{\chi}_{k} \chi_{k}-\xi \bar{\chi}_{k}^{\prime} \chi_{k}^{\prime}+\bar{\chi}_{k}^{\prime} \chi_{k+1}^{\prime}-\bar{\chi}_{k} \chi_{k+1}+\bar{\chi}_{k+1}^{\prime} \chi_{k}^{\prime}-\bar{\chi}_{k+1} \chi_{k}\right] \tag{7.184}
\end{align*}
$$

by finishing the Grassmann integration at the "time moment" $n-1$. In order to distinguish between the real density matrix and the fictitious "bosonic" and "fermionic density matrices" (7.184), the latter is indicated by a bar. The "fermionic density matrix" is a polynomial consisting of 16 terms where each of four coordinates $\left(\bar{\chi}_{n}, \chi_{n}, \bar{\chi}_{n}^{\prime}, \chi_{n}^{\prime}\right)$ occurs either in the first or in zero power, that is

$$
\begin{align*}
& \bar{\rho}_{n}\left(\bar{\chi}, \chi, \bar{\chi}^{\prime}, \chi^{\prime}\right)=\bar{\rho}_{0000}(n)+\chi^{\prime} \bar{\rho}_{0001}(n) \\
& +\bar{\chi}^{\prime} \bar{\rho}_{0010}(n)+\bar{\chi}^{\prime} \chi^{\prime} \bar{\rho}_{0011}(n)+\chi \bar{\rho}_{0100}(n) \\
& +\chi \chi^{\prime} \bar{\rho}_{0101}(n)+\chi \bar{\chi}^{\prime} \bar{\rho}_{0110}(n)+\chi \bar{\chi}^{\prime} \chi^{\prime} \bar{\rho}_{0111}(n) \\
& +\overline{\chi \rho}_{1000}(n)+\bar{\chi} \chi^{\prime} \bar{\rho}_{1001}(n) \\
& +\overline{\chi \chi}^{\prime} \bar{\rho}_{1010}(n)+\overline{\chi \chi}^{\prime} \chi^{\prime} \bar{\rho}_{1011}(n)+\bar{\chi} \chi \bar{\rho}_{1100}(n) \\
& +\bar{\chi} \chi \chi^{\prime} \bar{\rho}_{1101}(n)+\bar{\chi} \chi \bar{\chi}^{\prime} \bar{\rho}_{1110}(n)+\bar{\chi} \chi \bar{\chi}^{\prime} \chi^{\prime} \bar{\rho}_{1111}(n) \tag{7.185}
\end{align*}
$$

With the help of the leftmost line vector ( $1, \chi, \bar{\chi}, \bar{\chi} \chi$,$) and the corresponding$ rightmost column vector this can also be written as a bi-linear form

$$
\left.\begin{array}{rl} 
& \bar{\rho}_{n}\left(\bar{\chi}, \chi, \bar{\chi}^{\prime}, \chi^{\prime}\right)=(1, \chi, \bar{\chi}, \bar{\chi} \chi,) \\
\times & \left(\begin{array}{ccc}
\bar{\rho}_{0000} & \bar{\rho}_{0001} & \bar{\rho}_{0010} \\
\bar{\rho}_{0100} & \bar{\rho}_{0101} \\
\bar{\rho}_{1000} & \bar{\rho}_{1001} & \bar{\rho}_{1010} \\
\bar{\rho}_{0111} \\
\bar{\rho}_{1100} & \bar{\rho}_{1101} & \bar{\rho}_{1110}
\end{array} \bar{\rho}_{1111}\right.
\end{array}\right)\left(\begin{array}{c}
1  \tag{7.186}\\
\chi^{\prime} \\
\bar{\chi}^{\prime} \\
\bar{\chi}^{\prime} \chi^{\prime}
\end{array}\right) .
$$

where the coefficients forming a $4 \times 4$ matrix and depending on $n$ encode all of the fermionic information. Later, we will consider the evolution of a multiplet consisting of both "fermions" and "bosons", for which all matrix elements
are functions of bosonic variables $(x, \bar{x} ; y, \bar{y})$, or in short $\left(z, z^{\prime}\right)$. We attribute variables with primes to "antiparticles"; $z^{\prime}$ to "antibosons", and $\chi^{\prime}$ and $\bar{\chi}^{\prime}$ to "antifermions". We assume that "bosonic" coordinates are implicit in the representation (7.186).

We determine the recurrence relation for this matrix propagating it in "time" to $n+1$. The result of such a propagation is known: it should be equivalent to the transformation given by a $2 \times 2$ matrix in (7.168) for $E_{n}=$ 0 , which has to be applied to each of the participating fermions, that is separately for each index $(0,1)$, of the matrix elements. We consider this in more detail and show how this transformation appears in the matrix notation (7.186). We perform integration over all but the last four anticommuting variables and find the recurrence relations

$$
\begin{align*}
\bar{\rho}_{n}\left(\bar{\chi}, \chi, \bar{\chi}^{\prime}, \chi^{\prime}\right)= & \int d \chi_{1} d \bar{\chi}_{1} d \chi_{1}^{\prime} d \bar{\chi}_{1}^{\prime} \bar{\rho}_{n-1}\left(\bar{\chi}_{1}, \chi_{1}, \bar{\chi}_{1}^{\prime}, \chi_{1}^{\prime}\right)  \tag{7.187}\\
& \exp \left(\varepsilon \bar{\chi} \chi-\xi \bar{\chi}^{\prime} \chi^{\prime}+\bar{\chi}^{\prime} \chi_{1}^{\prime}-\bar{\chi} \chi_{1}+\bar{\chi}_{1}^{\prime} \chi^{\prime}-\bar{\chi}_{1} \chi\right)
\end{align*}
$$

for the functions of the Grassmann variables.
Let us first consider the result of integration over one variable, say $d \bar{\chi}_{1}^{\prime}$, which involves only the last column vector of the quadratic form (7.186). We take into account the anticommutation relations (7.178) among the differentials $d \chi_{1} d \bar{\chi}_{1} d \chi_{1}^{\prime} d \bar{\chi}_{1}^{\prime}=-d \bar{\chi}_{1}^{\prime} d \chi_{1}^{\prime} d \chi_{1} d \bar{\chi}_{1}$, allow for the relation $\exp \left(\bar{\chi}_{1}^{\prime} \chi^{\prime}\right)=$ $\left(1+\bar{\chi}_{1}^{\prime} \chi^{\prime}\right)$, and perform the integration following the rules (7.177). For the last vector of (7.186) this yields

$$
\left(\begin{array}{c}
1  \tag{7.188}\\
\chi_{1}^{\prime} \\
\bar{\chi}_{1}^{\prime} \\
\bar{\chi}_{1}^{\prime} \chi_{1}^{\prime}
\end{array}\right) \stackrel{d \bar{\chi}_{1}^{\prime}}{\rightarrow}\left(\begin{array}{c}
-\int\left(1+\bar{\chi}_{1}^{\prime} \chi^{\prime}\right) d \bar{\chi}_{1}^{\prime} \\
-\int \chi_{1}^{\prime}\left(1+\bar{\chi}_{1}^{\prime} \chi^{\prime}\right) d \bar{\chi}_{1}^{\prime} \\
-\int \bar{\chi}_{1}^{\prime}\left(1+\bar{\chi}_{1}^{\prime} \chi^{\prime}\right) d \bar{\chi}_{1}^{\prime} \\
-\int \bar{\chi}_{1}^{\prime} \chi_{1}^{\prime}\left(1+\bar{\chi}_{1}^{\prime} \chi^{\prime}\right) d \bar{\chi}_{1}^{\prime}
\end{array}\right)=\left(\begin{array}{c}
\chi^{\prime} \\
\chi_{1}^{\prime} \chi^{\prime} \\
-1 \\
\chi_{1}^{\prime}
\end{array}\right) .
$$

Continuing by integration over $d \chi_{1}^{\prime}$ allowing for the relation $\exp \left(\bar{\chi}^{\prime} \chi_{1}^{\prime}\right)=$ $\left(1+\bar{\chi}^{\prime} \chi_{1}^{\prime}\right)$ we obtain

$$
\left(\begin{array}{c}
\chi^{\prime}  \tag{7.189}\\
\chi_{1}^{\prime} \chi^{\prime} \\
-1 \\
\chi_{1}^{\prime}
\end{array}\right) \stackrel{d \chi_{1}^{\prime}}{\rightarrow}\left(\begin{array}{c}
\int \chi^{\prime}\left(1+\bar{\chi}^{\prime} \chi_{1}^{\prime}\right) d \chi_{1}^{\prime} \\
\int \chi_{1}^{\prime} \chi^{\prime}\left(1+\bar{\chi}^{\prime} \chi_{1}^{\prime}\right) d \chi_{1}^{\prime} \\
-\int \chi_{1}^{\prime}\left(1+\bar{\chi}^{\prime} \chi_{1}^{\prime}\right) d \chi_{1}^{\prime} \\
\int \chi^{\prime}\left(1+\bar{\chi}_{1}^{\prime} \chi_{1}^{\prime}\right) d \chi_{1}^{\prime}
\end{array}\right)=\left(\begin{array}{c}
-\bar{\chi}^{\prime} \chi^{\prime} \\
-\chi^{\prime} \\
-\bar{\chi}^{\prime} \\
1
\end{array}\right) .
$$

In other words, replacement of the variables with subscript 1 by variables without the subscript, is equivalent to multiplication of the vector column by a transformation matrix

$$
\left(\begin{array}{c}
1  \tag{7.190}\\
\chi_{1}^{\prime} \\
\bar{\chi}_{1}^{\prime} \\
\bar{\chi}_{1}^{\prime} \chi_{1}^{\prime}
\end{array}\right) \stackrel{d \bar{\chi}_{1}^{\prime} d \chi_{1}^{\prime}}{ }\left(\begin{array}{cccc}
0 & 0 & 0 & -1 \\
0 & -1 & 0 & 0 \\
0 & 0 & -1 & 0 \\
1 & 0 & 0 & 0
\end{array}\right)\left(\begin{array}{c}
1 \\
\chi^{\prime} \\
\bar{\chi}^{\prime} \\
\bar{\chi}^{\prime} \chi^{\prime}
\end{array}\right) .
$$

Let us allow now for multiplication by $\exp \left(-\xi \bar{\chi}^{\prime} \chi^{\prime}\right)=1-\xi \bar{\chi}^{\prime} \chi^{\prime}$ which yields

$$
\left(\begin{array}{c}
1  \tag{7.191}\\
\chi^{\prime} \\
\bar{\chi}^{\prime} \\
\bar{\chi}^{\prime} \chi^{\prime}
\end{array}\right) \stackrel{\exp \left(-\xi \bar{\chi}^{\prime} \chi^{\prime}\right)}{\rightarrow}\left(\begin{array}{c}
1-\xi \bar{\chi}^{\prime} \chi^{\prime} \\
\chi^{\prime} \\
\bar{\chi}^{\prime} \\
\bar{\chi}^{\prime} \chi^{\prime}
\end{array}\right)
$$

and therefore modifies (7.190)

$$
\left(\begin{array}{c}
1  \tag{7.192}\\
\chi^{\prime} \\
\bar{\chi}^{\prime} \\
\bar{\chi}^{\prime} \chi^{\prime}
\end{array}\right)_{1} d \bar{\chi}_{1}^{\prime} d \chi_{1}^{\prime} \times \exp \left(-\xi \bar{\chi}^{\prime} \chi^{\prime}\right)\left(\begin{array}{cccc}
-\xi & 0 & 0 & -1 \\
0 & -1 & 0 & 0 \\
0 & 0 & -1 & 0 \\
1 & 0 & 0 & 0
\end{array}\right)\left(\begin{array}{c}
1 \\
\chi^{\prime} \\
\bar{\chi}^{\prime} \\
\bar{\chi}^{\prime} \chi^{\prime}
\end{array}\right)
$$

We can calculate the transformation of the line vectors analagously and obtain

$$
(1, \chi, \bar{\chi}, \bar{\chi} \chi,)_{1} \rightarrow(1 \chi, \bar{\chi}, \bar{\chi} \chi,)\left(\begin{array}{cccc}
0 & 0 & 0 & 1  \tag{7.193}\\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0 \\
-1 & 0 & 0 & \varepsilon
\end{array}\right)
$$

In other words, the matrix $\widehat{\bar{\rho}}$ of the coefficients $\bar{\rho}_{i j k l}$ of (7.186) experience the
 and (7.193), respectively.

Note that the fermionic part of the line and the column vectors that contain only linear terms $\chi^{\prime}, \bar{\chi}^{\prime}$ or $\chi, \bar{\chi}$ is not affected by the transformation over the period and does not change, apart from an unimportant variation of the signs. This allows us to concentrate only on the first and the last lines and columns of the matrices, discard the second and the third line and column and write

$$
\begin{align*}
& \binom{1}{\bar{\chi}^{\prime} \chi^{\prime}}_{1} \rightarrow\left(\begin{array}{cc}
-\xi, & -1 \\
1, & 0
\end{array}\right)\binom{1}{\bar{\chi}^{\prime} \chi^{\prime}} .  \tag{7.194}\\
& (1, \bar{\chi} \chi)_{1} \rightarrow(1, \bar{\chi} \chi)\left(\begin{array}{cc}
0, & 1 \\
-1, & \varepsilon
\end{array}\right) \tag{7.195}
\end{align*}
$$

which is equivalent to (7.168). In other words, combinations of "particles" containing an odd number of "fermions" or "antifermions" are not affected by the "time" evolution. Neither are they affected by the disorder. This circumstance is very important to note, since the operators $\bar{\chi}_{0}^{\prime} \bar{\chi}_{0}$ and $\chi_{n} \chi_{n}^{\prime}$ in front of the exponent in the integrand (7.182) do change the number of "fermions" from even to odd. This occurs at the "time moments" corresponding to the starting level 0 and the final level $n$ in our original problem of the population dynamics, and means that during this "interval of time", the propagation of the "fermions" is disentangled from that of the "bosons", where the latter propagate independently according to (7.166).

After discarding from (7.186) the lines and the columns corresponding to odd numbers of fermions and antifermions, the "density matrix" for fictitious particles adopts the form

$$
\begin{align*}
& \overline{\widetilde{\rho}}_{n}\left(\bar{\chi}, \chi, x, \bar{x} ; \bar{\chi}^{\prime}, \chi^{\prime}, y, \bar{y}\right)=(1, \bar{\chi} \chi) \hat{\rho}_{n}\left(z, z^{\prime}\right)\binom{1}{\bar{\chi}^{\prime} \chi^{\prime}} ; \\
& \widehat{\rho}_{n}\left(z, z^{\prime}\right)=\binom{\widetilde{\bar{\rho}}_{0000}\left(z, z^{\prime}, n\right), \widetilde{\bar{\rho}}_{0011}\left(z, z^{\prime}, n\right)}{\bar{\rho}_{1100}\left(z, z^{\prime}, n\right), \overline{\bar{\rho}}_{1111}\left(z, z^{\prime}, n\right)} \tag{7.196}
\end{align*}
$$

where the bosonic variables $(x, \bar{x} ; y, \bar{y})$ are written explicitly as complex numbers $\left(z, z^{\prime}\right)$, as in (7.136). At $\gamma=0$ according to (7.162), (7.194), and (7.195), the evolution of the "density matrix" over the period reads

$$
\begin{equation*}
\widehat{\rho}_{n+1}\left(z, z^{\prime}\right)=\widehat{U}_{l} \widehat{u}_{L} \hat{\rho}_{n}\left(z, z^{\prime}\right) \widehat{u}_{R} \widehat{U}_{r} \tag{7.197}
\end{equation*}
$$

where

$$
\begin{align*}
\widehat{u}_{R} & =\left(\begin{array}{cc}
-\xi & -1 \\
1 & 0
\end{array}\right) ; \widehat{u}_{L}=\left(\begin{array}{cc}
0 & 1 \\
-1 & \varepsilon
\end{array}\right) \\
\widehat{U}_{r} & =\exp \left\{i \frac{1}{4}|\widehat{\mid}|_{z^{\prime}}^{2}\right\} \exp \left\{i(2-\xi)\left|z^{\prime}\right|^{2}\right\} \\
\widehat{U}_{l} & =\exp \left\{-i(2-\varepsilon)|z|^{2}\right\} \exp \left\{-i \frac{1}{4}|\widehat{p}|_{z}^{2}\right\} . \tag{7.198}
\end{align*}
$$

### 7.6.2 Mapping over a Period and the Ensemble Average

For $\gamma \neq 0$, equation (7.183) results from the ensemble averaging of the expression

$$
\begin{align*}
\rho_{n}(t)= & \int_{-\infty-i \nu}^{\infty-i \nu} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \xi d \varepsilon}{4 \pi^{2}} \prod_{k=-\infty}^{\infty}\left\{\int_{-\infty}^{\infty} \frac{d z_{k}^{*} d z_{k} d z_{k}^{\prime *} d z_{k}^{\prime}}{4 \pi^{2}}\right\} \\
& \prod_{k=-\infty}^{\infty}\left\{\int d \chi_{k} d \bar{\chi}_{k} d \chi_{k}^{\prime} d \bar{\chi}_{k}^{\prime}\right\} \bar{\chi}_{0}^{\prime} \chi_{n}^{\prime} \bar{\chi}_{0} \chi_{n} e^{-i(\varepsilon-\xi) t+S} ; \\
S= & \sum_{k=-\infty}^{\infty}\left[\left(\varepsilon-\gamma_{k}\right) \bar{\chi}_{k} \chi_{k}-\left(\xi-\gamma_{k}\right) \bar{\chi}_{k}^{\prime} \chi_{k}^{\prime}+i\left(\varepsilon-\gamma_{k}\right)\left|z_{k}\right|^{2}\right. \\
& \left.-i\left(\xi-\gamma_{k}\right)\left|z_{k}^{\prime}\right|^{2}+i z_{k}^{\prime} z_{k \pm 1}^{\prime *}-i z_{k} z_{k \pm 1}^{*}+\bar{\chi}_{k}^{\prime} \chi_{k \pm 1}^{\prime}-\bar{\chi}_{k} \chi_{k \pm 1}\right], \tag{7.199}
\end{align*}
$$

with the distribution function (7.130) for $E_{n}$. Here, as earlier, $z_{k}=x_{k}+i \bar{x}_{k}$, $z_{k}^{\prime}=y_{k}+i \bar{y}_{k}$. This suggests the consideration of the recurrence relations for fictitious "density matrices" $\bar{\rho}_{n}(7.186)$ for each individual realization of the disorder

$$
\begin{align*}
& \bar{\rho}_{n}\left(\bar{\chi}_{n}, \chi_{n}, x_{n}, \bar{x}_{n} ; \bar{\chi}_{n}^{\prime}, \chi_{n}^{\prime}, y_{n}, \bar{y}_{n}\right)=\prod_{k=-\infty}^{n-1}\left\{\int_{-\infty}^{\infty} \frac{d x_{k} d y_{k} d \bar{x}_{k} d \bar{y}_{k}}{\pi^{2}}\right\} \\
& \times \prod_{k=-\infty}^{n-1}\left\{\int d \chi_{k} d \bar{\chi}_{k} d \chi_{k}^{\prime} d \bar{\chi}_{k}^{\prime}\right\} \exp \sum_{k=-\infty}^{n}\left[2 i x_{k} x_{k-1}+2 i \bar{x}_{k} \bar{x}_{k-1}-2 i y_{k} y_{k-1}\right. \\
& -2 i \bar{y}_{k} \bar{y}_{k-1}+\bar{\chi}_{k}^{\prime} \chi_{k-1}^{\prime}-\bar{\chi}_{k} \chi_{k-1}+\bar{\chi}_{k-1}^{\prime} \chi_{k}^{\prime}-\bar{\chi}_{k-1} \chi_{k}+\left(\varepsilon-E_{k}\right) \bar{\chi}_{k} \chi_{k} \\
& \left.-\left(\xi-E_{k}\right) \bar{\chi}_{k}^{\prime} \chi_{k}^{\prime}+i\left(\varepsilon-E_{k}\right)\left(x_{k}^{2}+\bar{x}_{k}^{2}\right)-i\left(\xi-E_{k}\right)\left(y_{k}^{2}+\bar{y}_{k}^{2}\right)\right] \tag{7.200}
\end{align*}
$$

where the integrations are performed over all variables with indices less than $n$. Since the result of such a multistep Gaussian integration is a Gaussian function, one can assume that the "density matrix" has the form

$$
\begin{align*}
\bar{\rho}_{n-1}= & \exp \left[i M_{11}^{(n-1)}\left|z_{n-1}\right|^{2}-i M_{22}^{(n-1)}\left|z_{n-1}^{\prime}\right|^{2}\right. \\
& \left.-M_{33}^{(n-1)} \bar{\chi}_{n-1} \chi_{n-1}+M_{44}^{(n-1)} \bar{\chi}_{n-1}^{\prime} \chi_{n-1}^{\prime}\right] \tag{7.201}
\end{align*}
$$

and after substitution of (7.201) into (7.200), followed by integration, we find a Gaussian function again for the next $n$

$$
\begin{equation*}
\bar{\rho}_{n}=\exp \left[i M_{11}^{(n)}\left|z_{n}\right|^{2}-i M_{22}^{(n)}\left|z_{n}^{\prime}\right|^{2}-M_{33}^{(n)} \bar{\chi}_{n} \chi_{n}+M_{44}^{(n)} \bar{\chi}_{n}^{\prime} \chi_{n}^{\prime}\right] \tag{7.202}
\end{equation*}
$$

where the diagonal matrices $\widehat{M}^{(n)}$ and $\widehat{M}^{(n-1)}$ are defined as

$$
\widehat{M}^{(n)}=\left(\begin{array}{cccc}
M_{11}^{(n)} & 0 & 0 & 0 \\
0 & M_{22}^{(n)} & 0 & 0 \\
0 & 0 & M_{33}^{(n)} & 0 \\
0 & 0 & 0 & M_{44}^{(n)}
\end{array}\right)
$$

and satisfy the recurrence relation

$$
\widehat{M}^{(n)}=-\left(\widehat{M}^{(n-1)}\right)^{-1}+\left(\begin{array}{cccc}
\varepsilon-E_{n} & 0 & 0 & 0  \tag{7.203}\\
0 & \xi-E_{n} & 0 & 0 \\
0 & 0 & \varepsilon-E_{n} & 0 \\
0 & 0 & 0 & \xi-E_{n}
\end{array}\right)
$$

This is a direct analog of mapping (7.110). One sees that in this construction, $M_{11}^{(n)}$ always coincides with $M_{33}^{(n)}$, whereas $M_{22}^{(n)}$ coincides with $M_{44}^{(n)}$. We have already encountered such a situation considering the tunneling transparency.

Successive application of the mapping yields the matrix $\widehat{M}^{(n)}$ as a function of the initial condition $\widehat{M}^{(0)}$ and all of the string of intermediate energies $E_{k}$ with $0<k<n$. Averaging over the ensemble of random energies $E_{k}$ therefore generates an ensemble of random matrices $\widehat{M}^{(n)}$ with a certain distribution function $g\left(M_{11}^{(n)}, M_{22}^{(n)}, \varepsilon, \xi\right)$. For an infinitely long chain of levels,
this distribution function loses the memory of the initial condition $\widehat{M}^{(0)}$ and becomes "stationary", not depending on the level number $n$, that is on the fictitious "time". Hence for the ensemble average "density matrix" one finds

$$
\begin{align*}
& \bar{\rho}_{n}\left(\bar{\chi}, \chi, z ; \bar{\chi}^{\prime}, \chi^{\prime}, z^{\prime}\right)=\int d M_{11} d M_{22} g\left(M_{11}, M_{22}, \varepsilon, \xi\right) \\
& \exp \left[i M_{11}|z|^{2}-i M_{22}\left|z^{\prime}\right|^{2}-M_{11} \bar{\chi} \chi+M_{22} \bar{\chi}^{\prime} \chi^{\prime}\right] \tag{7.204}
\end{align*}
$$

where the "stationary" distribution function $g\left(M_{11}^{(n)}, M_{22}^{(n)}, \varepsilon, \xi\right)$ of matrices satisfies the equation

$$
\begin{align*}
g\left(M_{11}, M_{22}, \varepsilon, \xi\right)= & \int d E_{n} \frac{\tilde{g}\left(E_{n}\right)}{\left(\varepsilon-E_{n}-M_{11}\right)^{2}}  \tag{7.205}\\
& \frac{1}{\left(\xi-E_{n}-M_{22}^{2}\right)^{2}} g\left(\frac{1}{\varepsilon-E_{n}-M_{11}}, \frac{1}{\xi-E_{n}-M_{22}}, \varepsilon, \xi\right),
\end{align*}
$$

which follows from the recurrence relation (7.203) and the random energy distribution (7.130)

$$
\widetilde{g}\left(E_{n}\right)=\frac{1}{\sqrt{\pi} \gamma} \exp \frac{-E_{n}^{2}}{\gamma^{2}}
$$

Integration over $d M_{11} d M_{22}$ can be referred to as integration over a matrix $d \widehat{M}$, which for the particular case under consideration is reduced to integration over only two diagonal matrix elements. For the real numbers $\varepsilon, \xi, E_{n}$ the matrix elements $M_{11}$ and $M_{22}$ are real, but even a vanishing imaginary part of $\varepsilon$ and $\xi$ makes these matrix elements complex when mapping (7.203) is performed an infinite number of times. In the latter case the symbol $d M_{11} d M_{22}$ implies independent integration over the real and the imaginary parts of $M_{11}$ and $M_{22}$, for which one always has two constraints $\operatorname{Im} M_{11}>0>\operatorname{Im} M_{22}$. The distribution $g\left(M_{11}, M_{22}, \varepsilon, \xi\right)$ in this case is a real function of the complex variables $M_{11}, M_{22}$.

Note that the representation

$$
\begin{align*}
\widetilde{\rho}_{n}\left(z, z^{\prime}\right)= & \int d M_{11} d M_{22} \bar{g}_{n}\left(M_{11}, M_{22}, \varepsilon, \xi\right) \times \\
& \exp \left[i M_{11}|z|^{2}-i M_{22}\left|z^{\prime}\right|^{2}\right] \tag{7.206}
\end{align*}
$$

similar to (7.204) can be given for the purely "bosonic" density matrix (7.161) although with a slightly different recurrence relation

$$
\begin{align*}
\bar{g}_{n+1}\left(M_{11}, M_{22}, \varepsilon, \xi\right)= & \int d E_{n} \frac{\widetilde{g}\left(E_{n}\right)}{\varepsilon-E_{n}-M_{11}} \\
& \frac{1}{\xi-E_{n}-M_{22}} \bar{g}_{n}\left(\frac{1}{\varepsilon-E_{n}-M_{11}}, \frac{1}{\xi-E_{n}-M_{22}}, \varepsilon, \xi\right), \tag{7.207}
\end{align*}
$$

for the distribution function. The difference originates in the fact that integration over bosonic variables $z, z^{\prime}$ yields denominators $M_{11}^{\prime} M_{22}^{\prime}$ that are not compensated by the same factors in numerators resulting from Grassmann integrations. Due to this difference, the normalization of the distribution function (7.207) is not conserved, but exponentially decreases over the course of iterations, where the function itself does not have a non-zero asymptotic profile. The decrease rate is given by the modulus of the smallest eigenvalue $e^{\lambda(\varepsilon, \xi)}$ of the integral operator (7.207), and we denote the corresponding eigenfunction by $\bar{g}_{\text {min }}\left(M_{11}^{\prime}, M_{22}^{\prime}, \varepsilon, \xi\right)$.


Fig. 7.10. Distribution of the points $M_{11}$ and $M_{22}$ in the upper and lower parts of the complex plane after a large number of mappings (7.203) with randomly distributed $E_{n}$. (a) A Lorentzian stationary distribution is attained close to the real axis for the case $\operatorname{Im} \zeta / \gamma^{2} \ll 1$, whereas in the opposite limit (b) one sees a narrow Gaussian distribution around the attracting stable points. The corresponding distribution functions $g\left(M_{11}\right)$ are shown in (c) and (d) respectively.

Equation (7.205) implies that a stationary distribution $g\left(M_{11}, M_{22}, \varepsilon, \xi\right)$ is attained as a result of the interplay between the compressing mapping (7.203) at $E_{n}=0$ and the random walk resulting from the irregular energy shifts $E_{n}$. Mapping in the complex planes $M_{11}$ and $M_{22}$ tends to compress the distribution to the stable stationary points shown for (7.110) in Fig. 7.7. For $\operatorname{Im}(\varepsilon-\xi)<0$ these points are located outside of the unit circles $\left|M_{11}\right|=1$ and $\left|M_{22}\right|=1$ at a distance $|\operatorname{Im} \zeta| / \sqrt{4-\eta^{2}}$, which can be found for small $\zeta$ by straightforward calculation of the positions of the stationary points. Here $\eta=$ $(\varepsilon+\xi) / 2$ and $\zeta=\varepsilon-\xi$, as usual. The compression factor $1-2|\operatorname{Im} \zeta| / \sqrt{4-\eta^{2}}$ narrows the distribution around the stationary points, whereas the random
values of $E_{n}$ result in a diffusive broadening of the distribution. Relative broadening of a Gaussian distribution $\exp \left[-a\left|M_{i i}-M_{s t}\right|^{2}\right]$ as a result of one additional random diffusion step of length $\gamma$ amounts to $\gamma^{2} a+1$, and therefore is completely compensated by the compression if $\gamma^{2} a=2|\operatorname{Im} \zeta| / \sqrt{4-\eta^{2}}$. Hence, for $|\operatorname{Im} \zeta| / \gamma^{2} \gg 1$ the stationary distribution is a gaussian

$$
\begin{align*}
g\left(M_{11}, M_{22}, \varepsilon, \xi\right)= & \exp \left[-2|\operatorname{Im} \zeta| \frac{\left|M_{11}-e^{i \arccos \eta / 2}\right|^{2}}{\gamma^{2} \sqrt{4-\eta^{2}}}\right] \\
& \exp \left[-2|\operatorname{Im} \zeta| \frac{\left|M_{22}-e^{-i \arccos \eta / 2}\right|^{2}}{\gamma^{2} \sqrt{4-\eta^{2}}}\right] \frac{\gamma^{4}\left(4-\eta^{2}\right)}{4 \pi^{2}|\operatorname{Im} \zeta|^{2}}, \tag{7.208}
\end{align*}
$$

where we have determined the positions of the stationary points

$$
M_{s t}=\exp ( \pm i \arccos \eta / 2)
$$

in the approximation $|\zeta| \ll \eta$ and have allowed for the normalization factor.
The inverse Fourier transformation is performed however along the contours, where the imaginary parts of both $\varepsilon$ and $\xi$ are small as compared to any physical parameter. Therefore the other extreme $|\operatorname{Im} \zeta| / \gamma^{2} \ll 1$ is important. In this limit, the distribution tends to its asymptotic form

$$
\begin{equation*}
g\left(M_{11}, M_{22}, \varepsilon, \xi\right)=\frac{\sqrt{\left(4-\varepsilon^{2}\right)\left(4-\xi^{2}\right)}}{\pi^{2}\left(2+2 M_{11}^{2}-2 M_{11} \varepsilon\right)\left(2+2 M_{22}^{2}-2 M_{22} \xi\right)} \tag{7.209}
\end{equation*}
$$

which is a solution of $(7.205)$ for $\gamma \rightarrow 0,(\operatorname{Im} \zeta) / \gamma \rightarrow 0$ in complete agreement with the distribution (7.107), (7.113). In Fig. 7.10 we show the results of a large number of sequential mappings (7.203) performed numerically, with randomly chosen sequences of $E_{n}$ for both the $|\operatorname{Im} \zeta| / \gamma^{2} \gg 1$ and $|\operatorname{Im} \zeta| / \gamma^{2} \ll$ 1 cases. One sees that the values of $M_{11}$ and $M_{22}$ obtained for different realizations of the disorder from the same initial conditions are scattered in the upper and lower parts of the complex plane, respectively. Two regimes of small and large ratio $|\operatorname{Im} \zeta| / \gamma^{2}$ corresponding to distributions (7.208) and (7.209) respectively, are clearly seen.

For the "bosonic density matrix" (7.207) the distribution is similar. For the case $|\operatorname{Im} \zeta| / \gamma^{2} \ll 1$ one finds

$$
\begin{align*}
\bar{g}_{n}(\widehat{M}, \varepsilon, \xi) & =\bar{g}_{\min }(\widehat{M}, \varepsilon, \xi) e^{-n \lambda(\varepsilon, \xi)}, \\
\bar{g}_{\min }(\widehat{M}, \varepsilon, \xi) & =\frac{4}{\pi^{2}\left(\varepsilon+i \sqrt{4-\varepsilon^{2}}-2 M_{11}\right)\left(\xi-i \sqrt{4-\xi^{2}}-2 M_{22}\right)},  \tag{7.210}\\
\exp [\lambda(\varepsilon, \xi)] & =\left(\frac{\varepsilon}{2}+i \sqrt{1-\frac{\varepsilon^{2}}{4}}\right)\left(\frac{\xi}{2}-i \sqrt{1-\frac{\xi^{2}}{4}}\right)=e^{i \arccos \varepsilon / 2-i \arccos \xi / 2},
\end{align*}
$$

which satisfys the recurrence equation (7.207) at $\gamma=0$. We note that substitution of (7.209), (7.210) into (7.204), (7.206), respectively, yields the
"density matrices", which do not change as a result of the transformations (7.162), (7.198), thus remaining constants in the interaction representation.

Allowance for the finite width $\gamma$ of the distribution of $E_{n}$ can be considered as a perturbation of (7.209), (7.210). For the distribution function $\bar{g}_{\text {min }}$ of (7.210) this yields a small $\left(\sim \gamma^{2}\right)$ correction, and also results in a correction for the eigenvalue parameter $\lambda$ of the same order of magnitude. The latter is of major importance, since this very correction leads to the localization effect and thereby changes the form of the asymptotic population distribution in the long-time limit. Corrections resulting from the difference between the identity operator applied to $\bar{g}_{\min }$ and a convolution with a sharp symmetric distribution (7.130) are proportional to the second derivative of this function, which should be projected back to $\bar{g}_{\text {min }}$, according to the standard prescription of perturbation theory. Performing the calculations for $\zeta \ll 1$ one finds $\lambda=i \zeta / \sqrt{4-\eta^{2}}-\gamma^{2} / 2\left(4-\eta^{2}\right)$ which coincides with the index of the exponent suggested by the Liouville equation (7.166) in the interaction representation.

### 7.6.3 Time Dependence of the Population Distribution

Now we can interpret the time evolution of the population distribution $\rho_{n}(t)$ of the relay (7.183) in terms of "density matrices" (7.204), (7.206) for fictitious particles evolving in fictitious time $n$. We denote the running value of this fictitious time by the index $k$. From $k=-\infty$ until the moment $k=0$, the " $n$ evolution" corresponds to the joint propagation of bosons and bi-fermions, and is given by their stationary "density matrix" (7.204). Note that in this case the matrix contains only the part $\overline{\widetilde{\rho}}_{n}$ given by (7.196) of the entire density matrix (7.186). At the point $k=0$, the Grassmann operators $\bar{\chi}_{0}^{\prime} \bar{\chi}_{0}$ entering the pre-exponential factor in the integrand of (7.183) annihilate one of the fermions, and as a result, the remaining fermion becomes decoupled from the bosons, propagating without change as we have seen in the context of (7.194). Note in this case that only the matrix element $\bar{\rho}_{1010}$ of (7.186) is different from zero and

$$
\tilde{\bar{\rho}}_{k}\left(\bar{\chi}, \chi, z ; \bar{\chi}^{\prime}, \chi^{\prime}, z^{\prime}\right)=\widetilde{\rho}_{k}\left(z, z^{\prime}\right) \bar{\chi}_{k}^{\prime} \bar{\chi}_{k}
$$

The bosons, in turn, start to experience decay, and the dominating contribution to their density matrix $\widetilde{\rho}_{k}\left(z, z^{\prime}\right)(7.206)$ comes from the distribution given by the eigenfunction $\widetilde{g}_{\min }\left(M_{11}^{\prime}, M_{22}^{\prime}, \varepsilon, \xi\right)$ with the smallest eigenvalue $e^{\lambda(\varepsilon, \xi)}$ corresponding to the slowest decay. At the point $k=n$, the Grassmann operators $\chi_{n}^{\prime} \chi_{n}$ restore the fermion, and the further propagation again corresponds to the stationary "density matrix" (7.204). We illustrate this in Fig. 7.11

The population $\rho_{n}(t)$ can therefore be written as overlap integrals

$$
\begin{align*}
\rho_{n}(t)= & \int_{-\infty-i \nu}^{\infty-i \nu} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \xi d \varepsilon}{4 \pi^{2}} e^{-i(\varepsilon-\xi) t} e^{n \lambda(\varepsilon, \xi)} \int \frac{d z^{*} d z d z^{\prime *} d z^{\prime}}{4 \pi^{2}} d \chi d \bar{\chi} d \chi^{\prime} d \bar{\chi}^{\prime} \\
& \bar{\rho}\left(\bar{\chi}, \chi, z ; \bar{\chi}^{\prime}, \chi^{\prime}, z^{\prime}\right) \bar{\chi}^{\prime} \bar{\chi} \widetilde{\rho}\left(z, z^{\prime}\right) \chi^{\prime} \chi \tag{7.211}
\end{align*}
$$



Fig. 7.11. Representation of the localization with the help of the Liouville equation for a decaying compound. The dependence of the population on $n$ can be interpreted as a Liouvillian evolution of a compound comprising "bosons" and "fermions" along a fictitious time axis $n$. The stationary compound being perturbed at the point $k=0$ becomes decomposed to "fermions" propagating without decay and "bosons" exponentially decaying with the course of time. In order to determine the population $\rho_{n}$ at the point $n$ one has to find an overlap of the decaying density matrix at the point $k=n$ with the density matrix of the stationary compound.
of the "density matrices" Eqs. $(7.204,7.206)$ of fictitious particles, where we have taken into account that these "density matrices" do not change with $n$, apart from attenuation by a factor $e^{n \lambda(\varepsilon, \xi)}$. After substitution of the right-hand sides of these equations and after performing integrations over "fermionic" and "bosonic" coordinates, one obtains another representation of the overlap integral,

$$
\begin{align*}
\rho_{n}(t)= & \int_{-\infty-i \nu}^{\infty-i \nu} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \xi d \varepsilon}{4 \pi^{2}} e^{-i(\varepsilon-\xi) t-n \lambda(\varepsilon, \xi)} \\
& \int d M_{11} d M_{22} g\left(M_{11}, M_{22}, \varepsilon, \xi\right) \bar{g}_{\min }\left(M_{11}, M_{22}, \varepsilon, \xi\right), \tag{7.212}
\end{align*}
$$

now in terms of the distribution functions (7.205) and (7.207). Note that substituting of the distribution functions (7.209), (7.210), which are exact in the limit $\gamma \rightarrow 0$, into this equation, followed by integration over the matrix elements $d M_{11}, d M_{22}, d M_{11}^{\prime}$, and $d M_{22}^{\prime}$ along the real axis and over $d \xi d \varepsilon$ one finds

$$
\begin{equation*}
\rho_{n}(t)=\left|J_{n}(2 t)\right|^{2}, \tag{7.213}
\end{equation*}
$$

in complete agreement with (7.11) for the relay of resonant levels. Evidently in this limit, one cannot see any trace of the disorder.

In order to take into account the effect of the disorder, one therefore has to allow for the deviation of the distribution functions entering (7.205) and (7.207) from their asymptotic forms given by (7.209) and (7.210). Let us assume that for both $g$ and $\bar{g}_{\text {min }}$ this deviation can be expressed in the form of a convolution of the asymptotic distribution with a narrow function of width vanishing in the limit $\gamma \rightarrow 0$. Although each of the functions $g$ and $\bar{g}_{\text {min }}$ may have an individual narrow counterpart, the population distribution is governed by the convolution $g_{\gamma}$ of these two counterparts. This modifies (7.212)

$$
\begin{align*}
\rho_{n}(t)= & \int_{-\infty-i \nu}^{\infty-i \nu} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \xi d \varepsilon}{4 \pi^{2}} e^{-i(\varepsilon-\xi) t-n \lambda(\varepsilon, \xi)} \int d M_{11} d M_{22} d u g_{\gamma}(u) \\
& g\left(M_{11}+u, M_{22}+u, \varepsilon, \xi\right) \bar{g}_{\min }\left(M_{11}-u, M_{22}-u, \varepsilon, \xi\right) \tag{7.214}
\end{align*}
$$

and integration over $\widehat{M}$, after substitution of (7.209), (7.210), results in

$$
\begin{equation*}
\rho_{n}(t)=\int_{-\infty-i \nu}^{\infty-i \nu} \int_{-\infty+i \nu}^{\infty+i \nu} \frac{d \xi d \varepsilon}{4 \pi^{2}} e^{-i(\varepsilon-\xi) t-n \lambda(\varepsilon, \xi)} \int d u \frac{g_{\gamma}(u)}{\left(-i \sqrt{4-\varepsilon^{2}}+u\right)\left(i \sqrt{4-\xi^{2}}+u\right)}, \tag{7.215}
\end{equation*}
$$

where $\lambda \simeq i(\arccos \varepsilon / 2-\arccos \zeta / 2)-\gamma^{2} / 2\left(4-\eta^{2}\right)$. Note that replacement of the distribution $g_{\gamma}(u)$ by the Dirac $\delta$-function, followed by integration over $\varepsilon$ and $\xi$, immediately yields a result similar to that of (7.159).

After the replacements $\varepsilon \rightarrow \cos \phi$ and $\xi \rightarrow \cos \varphi$ we take into account the contributions of the poles at the points $\phi=\arcsin \frac{u}{2 i}, \phi=\pi-\arcsin \frac{u}{2 i}, \varphi=$ $-\arcsin \frac{u}{2 i}$, and $\varphi=\pi+\arcsin \frac{u}{2 i}$ to the integrals over $\phi$ and $\varphi$ and find

$$
\begin{align*}
\rho_{n}(t)= & \int d u \frac{g_{\gamma}(u) u^{2}}{\cos ^{2} \arcsin u / 2 i} \exp \left[-2 i n \arcsin \frac{u}{2 i}-\frac{n \gamma^{2}}{u^{2}}\right] \\
& -\operatorname{Re} \int d u \frac{g_{\gamma}(u) u^{2}}{\cos ^{2} \arcsin u / 2 i} \exp \left[-4 i t \cos \arcsin \frac{u}{2 i}-i n \pi-\frac{n \gamma^{2}}{u^{2}}\right] . \tag{7.216}
\end{align*}
$$

In the limit $n \rightarrow \infty, t \rightarrow \infty$, this integral has the asymptotic form resulting from the contribution of the saddle points $u=\left(2 \gamma^{2}\right)^{1 / 3}, u=\left(\frac{n \gamma^{2}}{i t}\right)^{1 / 3}$

$$
\begin{align*}
\rho_{n}(t) \simeq & \frac{A}{\sqrt{n}} \exp \left[-3(\gamma / 2)^{2 / 3} n\right]-B \frac{n^{5 / 6}}{t^{4 / 3}}\left|g_{\gamma}\left[\left(\frac{n \gamma^{2}}{i t}\right)^{1 / 3}\right]\right| \\
& \exp \left[-\frac{3^{3 / 2} \gamma^{2 / 3}}{2} n^{1 / 3} t^{2 / 3}\right] \cos \left[\frac{3 \gamma^{2 / 3}}{2} n^{1 / 3} t^{2 / 3}-n \pi-\frac{2}{3} \pi+\theta_{\gamma}\right], \tag{7.217}
\end{align*}
$$

where

$$
\begin{align*}
A & \simeq \sqrt{\frac{4 \pi}{3}} g_{\gamma}\left(2^{1 / 3} \gamma^{2 / 3}\right) 2^{5 / 6} \gamma^{5 / 3} \\
B & =\sqrt{\frac{2 \pi}{3}} \gamma^{5 / 3} ; \quad \theta_{\gamma}=\arg g_{\gamma}\left[\left(\frac{n \gamma^{2}}{i t}\right)^{1 / 3}\right] \tag{7.218}
\end{align*}
$$

This means that an exponentially localized distribution of localization length $\sim \gamma^{-2 / 3} / 3$ is attained, while the deviation from this distribution has the asymptotic time dependence $\exp \left(-\right.$ const $\left.n^{1 / 3} t^{2 / 3}\right)$. We note that the particular form of the distribution function $g_{\gamma}$ does not affect the shape of the asymptotic population distribution, but only influences the pre-exponential factor of the time-dependent deviation from this distribution, which vanishes with the course of time.

## 8 Composite Complex Quantum Systems

We now consider level grouping in multilevel systems other than the onedimensional relay of isolated levels discussed in the previous chapter. The population dynamics of such systems can display various manifestations that are typical of the systems already considered, such as exponential decay at short times, quantum recurrences and revivals, cooperative behavior, and quantum localization. Such systems can also demonstrate diffusive behavior. Typically, a problem of this type does not correspond to any hidden symmetry resulting in specific commutation relations among different parts of the Hamiltonian that comprise a closed subalgebra, and the evolution occurs in the entire Hilbert space, being constrained only by energy conservation. We first consider several problems for the relay of multilevel bands, for which diffusive evolution is strongly affected by correlations between the energy position and the coupling of the levels, and then address a general question concerning quantum dynamics of the systems, that in the classical limit manifest random walk behavior. We relate the quantum localization phenomenon to the dimensionality of those random walks responsible for classical returns and to the quantum recurrences. We also consider spectral manifestations typical of a quantum system perturbed by a random matrix, when both the structure of the unperturbed Hamiltonian and the mean squared coupling by the random matrix play an important role. The spectral transformation which occurs as a result of such perturbation explains many aspects of the time evolution of complex quantum systems.

### 8.1 Relay of Multilevel Bands

Consider a number of statistically equivalent, or identical spectral bands arranged in a one-dimensional array where each of the bands is coupled to each of its two neighbors via an interaction. The unperturbed Hamiltonian of the $n$-th band is given by its eigenvalues $E_{n, m}$ enumerated by the index $m$, while the coupling matrices are $V_{n, m}^{n+1, m^{\prime}}$ and $V_{n, m}^{n-1, m^{\prime}}$. The Schrödinger equation for the system reads

$$
\begin{equation*}
i \dot{\psi}_{n, k}=E_{n, k} \psi_{n, k}+\sum_{m}\left(V_{n, k}^{n+1, m} \psi_{n+1, m}+V_{n, k}^{n-1, m} \psi_{n-1, m}\right) \tag{8.1}
\end{equation*}
$$

and serves as a model for the physical problems corresponding to (2.4) and (2.56). In this section we consider three particular cases of (8.1). The first case relates to the situation where the total width of each of the bands is negligibly small as compared to the cooperative coupling, that is $\max \left(E_{n, k}-E_{n, k^{\prime}}\right) \ll N\left\langle V^{2}\right\rangle$ where $N$ is the total number of levels in each band and $\left\langle V^{2}\right\rangle$ is the mean squared interaction, as earlier. During a period of time shorter than $1 / \max \left(E_{n, k}-E_{n, k^{\prime}}\right)$ we can neglect the difference in level energies and consider the system as resonant. In the opposite extreme of broad bands, the total number $N$ of levels in each band plays an unimportant role, while the dynamics of the system is governed by the density $g$ of the band states. For time intervals shorter than the quantum recurrence time, the band populations satisfy the discrete diffusion equation or a master equation in the case of statistically non-equivalent bands.


Fig. 8.1. Relay of bands. (a) The bands can be considered as degenerate levels when the cooperative coupling $\left\langle V^{2}\right\rangle N \gg \max \left|\Delta_{n, m}-\Delta_{n, m^{\prime}}\right|$. (b) For the broad bands the population remains in the vicinities of resonances, whose total populations satisfy the master equation. (c) Correlation among the matrix elements $\left\langle V_{n, m}^{n+1, k} V_{n, m^{\prime}}^{n+1, k^{\prime}}\right\rangle \neq 0$ occurs when the rank of the interaction matrix is considerably smaller than its order. For rank 1 the correlation can be taken into account by introducing eigenstates of the coupling matrix, corresponding to non-vanishing eigenvalues. For the non-Hermitian matrix $V_{n, m}^{n+1, k}$ the right eigenvector $\left|C_{n}\right\rangle$ may not coincide with the left one $\left|C_{n}^{\prime}\right\rangle$. (d) When the right and the left eigenstates coincide the system resembles a relay of decaying levels.

Earlier, in Sec. 4.5, we have considered the influence of correlations among the coupling amplitudes and energy positions of the band levels on the population dynamics of two-band systems. The same influence exists in multiband problems, where the correlations are responsible for the initial stage of the evolution, which resembles the ballistic evolution of a classical particle pre-
ceding its diffusion in a disordered medium. As for the two-band systems, the correlations can manifest themselves in the ranks of the coupling matrices $V_{n, k}^{n \pm 1, m}$, that can be considerably smaller as compared to the order of the matrices. In this case one can choose the eigenstates of the coupling matrices corresponding to non-zero eigenvalues as basis states, appending this incomplete basis by orthogonal combination of the eigenstates of unperturbed band Hamiltonians, in the manner shown in Fig. 4.24b) for two-band systems. Four particular cases considered in this section are sketched in Fig. 8.1.

For the initial condition we always take $\psi_{n, k}(t=0)=\delta_{n}^{0} \delta_{k}^{0}$, since for any other Kronekker $\delta$-like initial condition the consideration is similar, whereas in general, initial conditions can be given as a linear combination of the $\delta$-like ones. We choose the energy position $E_{0,0}$ of the level $n=0, k=0$ as the energy reference point and denote $E_{n, k}-E_{0,0}=\Delta_{n, k}$. Fourier transformation in time of (8.1) therefore adopts the form
$\varepsilon \psi_{n, k}(\varepsilon)=\Delta_{n, k} \psi_{n, k}+\sum_{m}\left[V_{n, k}^{n+1, m} \psi_{n+1, m}(\varepsilon)+V_{n, k}^{n-1, m} \psi_{n-1, m}(\varepsilon)\right]+i \delta_{n}^{0} \delta_{k}^{0}$,
where, as earlier, we denote by the same symbol $\psi_{n, k}$ the time-dependent amplitude $\psi_{n, k}(t)$ and its Fourier transform $\psi_{n, k}(\varepsilon)$, since the argument clearly specifies which one of these two objects we have in mind. This equation is a generalization of (4.5) for the case of many bands.


Fig. 8.2.

Green's function, corresponding to (8.2), can be represented as a perturbation series (4.13) for the resolvent, in complete analogy with the twoband systems with random matrix elements, and can be depicted as diagrams Fig. 4.4, which after the ensemble average for the case $\left\langle V_{n, k}^{n+1, m} V_{n, k^{\prime}}^{n+1, m^{\prime}}\right\rangle=$ $\left\langle V^{2}\right\rangle \delta_{m}^{m^{\prime}} \delta_{k}^{k^{\prime}}$, adopt a tree-like structure (Fig. 4.7). This structure, resulting in the renormalized factors $X_{n}(\varepsilon)$ satisfying the equations

$$
\begin{align*}
X_{n, k}(\varepsilon) & =\frac{1}{\varepsilon-\Delta_{n, k}-\sum_{k}\left\langle V^{2}\right\rangle X_{n+1, k}(\varepsilon)-\sum_{k}\left\langle V^{2}\right\rangle X_{n-1, k}(\varepsilon)} \\
X_{n, k}(\xi) & =\frac{1}{\xi-\Delta_{n, k}-\sum_{k}\left\langle V^{2}\right\rangle X_{n+1, k}(\xi)-\sum_{k}\left\langle V^{2}\right\rangle X_{n-1, k}(\xi)} \tag{8.3}
\end{align*}
$$

by analogy to (4.20), (4.21), is shown in Fig. 4.8. This is also valid for the perturbation series for ensemble averaged populations, shown in Figs. 4.10
and 4.11. The only difference is in the fact that now instead of two alternating factors $X_{m}(\varepsilon) X_{m}(\xi)$ and $X_{n}(\varepsilon) X_{n}(\xi)$, we have a sequence of factors $X_{n, k}(\varepsilon) X_{n, k}(\xi)$ where the running indices $n$ for two neighboring factors are different by $\pm 1$, as shown in Fig. 8.2.

By analogy to (4.27) and Fig. 4.14 one can write the equation for the population

$$
\begin{align*}
\rho_{n}(\varepsilon, \xi)= & \rho_{n-1}(\varepsilon, \xi) \sum_{k}\left\langle V^{2}\right\rangle X_{n, k}(\varepsilon) X_{n, k}(\xi) \\
& +\rho_{n+1}(\varepsilon, \xi) \sum_{k}\left\langle V^{2}\right\rangle X_{n, k}(\varepsilon) X_{n, k}(\xi)+\delta_{n}^{0} \delta_{k}^{0} X_{0,0}(\varepsilon) X_{0,0}(\xi), \tag{8.4}
\end{align*}
$$

which corresponds to the diagram shown in Fig. 8.3.


Fig. 8.3.

### 8.1.1 Degenerate Bands With Random Coupling

For an infinite relay of degenerate bands shown in Fig. 8.1 a) the factors $X_{n, k}(\varepsilon)$ and $X_{n, k}(\xi)$ in (8.3) do not depend on the index $k$ thus satisfying the equations

$$
\begin{align*}
& X_{n}(\varepsilon)=\left[\varepsilon-\left\langle V^{2}\right\rangle N X_{n+1}(\varepsilon)-\left\langle V^{2}\right\rangle N X_{n-1}(\varepsilon)\right]^{-1} \\
& X_{n}(\xi)=\left[\xi-\left\langle V^{2}\right\rangle N X_{n+1}(\xi)-\left\langle V^{2}\right\rangle N X_{n-1}(\xi)\right]^{-1} \tag{8.5}
\end{align*}
$$

For the case of statistically equivalent bands, by analogy to (4.55) they have $n$-independent solutions

$$
\begin{align*}
& X_{n, k}(\varepsilon)=X(\varepsilon)=\frac{2}{\varepsilon-\sqrt{\varepsilon^{2}-8\left\langle V^{2}\right\rangle N}}, \\
& X_{n, k}(\xi)=X(\xi)=\frac{2}{\xi+\sqrt{\xi^{2}-8\left\langle V^{2}\right\rangle N}}, \tag{8.6}
\end{align*}
$$

different from (4.55) only by a factor of 2 in front of the number of states $N$ in the degenerate bands, which allows for the presence of two neighbors for each band.

There are many ways to solve (8.4), with the factors $X_{n, k}(\varepsilon)$ and $X_{n, k}(\xi)$ given by (8.6), that do not depend on indices. We choose one, based on the calculation of the trajectory sum for the diagrams of the type shown in Fig. 8.2. For this purpose, we note that the total number $T_{l, n}$ of trajectories
of length $l$ that start at the band number 0 and finish at the band number $n$, is given by the binomial coefficient

$$
\begin{equation*}
T_{l, n}=C_{l}^{(n+l) / 2}=\frac{l!}{\left(\frac{l+n}{2}\right)!\left(\frac{l-n}{2}\right)!}, \tag{8.7}
\end{equation*}
$$

as one can conclude from the sketch shown in Fig. 8.4 (a). Each trajectory



Fig. 8.4. (a) A number of random walk trajectories of length $l$ (dotted arrows) that result in the displacement $n$, given by the number of choices of $k+n$ steps upward and $k$ steps downward from the total number of $l=n+2 k$ steps. (b) Integration contours (dashed lines) for calculation of the generating function $\sum_{n=0}^{\infty} T_{l, n} x^{n}$ corresponding to these numbers. Poles of the integrand (8.11) are shown by circles.
of length $l$ brings a contribution $[X(\varepsilon) X(\xi)]^{l+1}\left(\left\langle V^{2}\right\rangle N\right)^{l}$ to the trajectory sum for populations, and therefore the net sum of all trajectories reads

$$
\begin{equation*}
\rho_{n}(\varepsilon, \xi)=\sum_{l=n}^{\infty} T_{l, n}[X(\varepsilon) X(\xi)]^{l+1}\left(\left\langle V^{2}\right\rangle N\right)^{l} \tag{8.8}
\end{equation*}
$$

The sum in (8.8) can be found with the help of the generating function

$$
\begin{equation*}
\Phi_{n}(x)=\sum_{l=n}^{\infty} \frac{l!}{\left(\frac{l+n}{2}\right)!\left(\frac{l-n}{2}\right)!} x^{l}=\sum_{l=0}^{\infty} \frac{l!}{\left(\frac{l+n}{2}\right)!\left(\frac{l-n}{2}\right)!} x^{l}, \tag{8.9}
\end{equation*}
$$

which can be extended on the integer $l<n$ since in this case, the second factorial in the denominator takes an infinite value. We calculate the numbers $T_{l, n}$ of (8.7) with the help of the integral representation

$$
\begin{equation*}
T_{l, n}=C_{l}^{(n+l) / 2}=\frac{1}{l+1} \frac{1}{B\left(\frac{l+n}{2}+1, \frac{l-n}{2}+1\right)}=\frac{2^{l}}{2 \pi} \int_{-\pi}^{\pi} d t \cos (n t) \cos ^{l} t \tag{8.10}
\end{equation*}
$$

for the $B$-function, which is valid for $l$ and $n$ of the same parity. Substitution of (8.10) into (8.9), followed by summation of the geometric series yields

$$
\begin{align*}
\Phi_{n}(x) & =\sum_{l=0}^{\infty} \frac{2^{l}}{2 \pi} \int_{0}^{2 \pi} d t \cos (n t) \cos ^{l} t x^{l} \\
& =\frac{1}{2 \pi} \int_{0}^{2 \pi} d t \frac{\cos (n t)}{1-2 x \cos t}=\frac{(2 x)^{n}}{2\left(1+\sqrt{1-4 x^{2}}\right)^{n} \sqrt{1-4 x^{2}}}, \tag{8.11}
\end{align*}
$$

where in the last line after the replacement $e^{i t} \rightarrow y$, the integration contour becomes a circle of unit radius around the point $y=0$ shown in Fig. 8.4b), separating the singularities of the integrand. For $n \geq 0$ integration along this contour gives the result directly, while for $n \leq 0$ the sign has to be changed.

Substituting $\varepsilon \rightarrow \sqrt{8\left\langle V^{2}\right\rangle N} \cos \phi, \xi \rightarrow \sqrt{8\left\langle V^{2}\right\rangle N} \cos \varphi$ into (8.6) we find the explicit expressions for the factors

$$
\begin{align*}
& 1 / X(\varepsilon)=\sqrt{2\left\langle V^{2}\right\rangle N} e^{-i \phi}, \\
& 1 / X(\xi)=\sqrt{2\left\langle V^{2}\right\rangle N} e^{i \varphi} \tag{8.12}
\end{align*}
$$

and making the replacement

$$
\begin{equation*}
x \rightarrow X(\varepsilon) X(\xi)\left(\left\langle V^{2}\right\rangle N\right)=\frac{1}{2} e^{i \phi-i \varphi} \tag{8.13}
\end{equation*}
$$

in (8.11) we obtain

$$
\begin{equation*}
\rho_{n}(\phi, \varphi)=\frac{e^{i(n+1)(\phi-\varphi)}}{2\left\langle V^{2}\right\rangle N\left(1+\sqrt{1-e^{2 i(\phi-\varphi)}}\right)^{n} \sqrt{1-e^{2 i(\phi-\varphi)}}} \tag{8.14}
\end{equation*}
$$

for the populations (8.8). The inverse Fourier transform of $\rho_{n}(\varepsilon, \xi)$ in the coordinates $(\phi, \varphi)$ therefore reads

$$
\begin{equation*}
\rho_{n}(\tau)=\frac{1}{\pi^{2}} \int_{C_{2} C_{1}} d \phi d \varphi \frac{\sin \phi \sin \varphi e^{i(n+1)(\phi-\varphi)-i \tau(\cos \phi-\cos \varphi)}}{\left(1+\sqrt{1-e^{2 i(\phi-\varphi)}}\right)^{n} \sqrt{1-e^{2 i(\phi-\varphi)}}} \tag{8.15}
\end{equation*}
$$

where $\tau=t \sqrt{8\left\langle V^{2}\right\rangle N}$, and the integration contours for $\phi$ and $\varphi$ coincide with the contours $C_{1}$ and $C_{2}$, respectively, shown in Fig. 4.17(b).

For further simplification we introduce the variables

$$
\begin{align*}
v & =\frac{\phi+\varphi}{2} \\
u & =\varphi-\phi . \tag{8.16}
\end{align*}
$$

Due to the periodicity of the integrand (8.15) the integration contours for $u$ and $v$ are the same intervals $(-\pi, \pi)$, and merely shifted by $-i 0$ for $u$. Note that this shift ensures the convergence of the sum (8.9) for the generating function. We therefore arrive at

$$
\begin{align*}
\rho_{n}(\tau) & =\int_{-\pi}^{\pi} \int_{-\pi-i 0}^{\pi-i 0} d u d v \frac{[\cos u-\cos 2 v] e^{i 2 \tau \sin (u / 2) \sin v} e^{-i(n+1) u}}{2 \pi^{2}\left(1+\sqrt{1-e^{-2 i u}}\right)^{n} \sqrt{1-e^{-2 i u}}} \\
& =\int_{-\pi-i 0}^{\pi-i 0} d u \frac{\cos u J_{0}\left(2 \tau \sin \frac{u}{2}\right)-J_{2}\left(2 \tau \sin \frac{u}{2}\right)}{\pi\left(1+\sqrt{1-e^{-2 i u}}\right)^{n} \sqrt{1-e^{-2 i u}}} e^{-i(n+1) u} \tag{8.17}
\end{align*}
$$

where in the last line the integral representation (4.60) of the Bessel functions $J_{k}(x)$ has been employed.


Fig. 8.5. Population distribution in an infinite relay of degenerate levels coupled by random matrices. Initially all the population is in level $n=0$. The time $t$ is scaled as $\tau=\sqrt{8 N\left\langle V^{2}\right\rangle} t$. (a) Time dependence of the population for the first several levels. (For negative $n$ the populations are identical.) (b) Dependence of population on time and index $n$. (Nonlinear scale $\sqrt{\rho_{n}}$ for the population.) (c) Universal asymptotic profile $\Phi(Y)$ of the population distribution (solid line) depends only on the parameter $Y^{2}=n^{2} / \tau$, the same as that which enters the diffusion profile, $\exp \left(-Y^{2} / 2\right)$ (dash-dot line).

In the asymptotic of large $n$ and large $\tau$ the main contribution comes from the domains $u \sim 1 / n^{2}$ and $u-\pi \sim 1 / n^{2}$ where the denominator is not an exponentially small value. Replacing $\left(1+\sqrt{1-e^{-2 i u}}\right)^{n}$ in these domains by its approximate value $\exp n \sqrt{2 i u}$ we find

$$
\begin{align*}
\rho_{n}(\tau)= & \int_{-\infty}^{\infty} d u \frac{J_{0}(\tau u)-J_{2}(\tau u)}{\pi \sqrt{2 i u}} e^{-n \sqrt{2 i u}} \\
& +\int_{-\infty}^{\infty} d u \frac{J_{0}(2 \tau)+J_{2}(2 \tau)}{\pi \sqrt{2 i u}} e^{-i n \pi} e^{-n \sqrt{2 i u}} \tag{8.18}
\end{align*}
$$

where the first and the second integral correspond to the first and the second domain, respectively, and the integration limits for these exponentially converging integrals are extended to infinity. For large $n$ and $\tau$ the second integral scales as $1 / n \sqrt{\tau}$ and can be ignored, as compared to the first integral, which scales as $1 / \sqrt{\tau}$. Replacing positive $u$ by $y^{2} / \tau$ and negative $u$ by $-y^{2} / \tau$ in the first integral we arrive at

$$
\begin{equation*}
\rho_{n}(\tau) \simeq \frac{2 \sqrt{2}}{\pi \sqrt{\tau}} \int_{0}^{\infty} d y\left[J_{0}\left(y^{2}\right)-J_{2}\left(y^{2}\right)\right] e^{-Y y} \cos (Y y+\pi / 4) \tag{8.19}
\end{equation*}
$$

where $Y=n / \sqrt{\tau}$. The integral on the right-hand side of (8.19) is a universal function $\sqrt{\pi} \Phi(Y) / 4$ of the parameter $Y$, which for large $Y$ is not very much different from a Gaussian function $\exp \left(-Y^{2} / 2\right)$, such that the population distribution (8.19) resembles a diffusive one.

In Fig. 8.5(a) we show the total population of the first few degenerate levels calculated with the help of the exact expression (8.17). One sees dying oscillations resembling those shown in Fig. 4.19(c) for two degenerate levels. Over the course of time, the population distribution becomes broader and its width increases as the square root of time, as one can see in Fig. 8.5(b), whereas the asymptotic profile $\Phi(Y)$ of the distribution is shown in Fig. 8.5(c).

The validity of the approach is limited in time by the assumption of the tree-like topological structure of the diagrams shown in Fig. 4.7. Indeed, for a one-dimensional random walk the mean width of the diffusive population distribution, and hence the mean displacement of the random trajectory is of the order of the square root of the random steps. In the context of the problem under consideration, this means that at the time $t$, a typical length of the trajectory is of the order of $\tau=t \sqrt{8\left\langle V^{2}\right\rangle N}$ and a typical width of the distribution is $n_{t} \sim \sqrt{\tau} \sim t^{1 / 2}\left(8\left\langle V^{2}\right\rangle N\right)^{1 / 4}$. This means that the trajectory returns $\tau / n_{t} \sim t^{1 / 2}\left(8\left\langle V^{2}\right\rangle N\right)^{1 / 4}$ times to the initially populated level, and this number of returns should remain much smaller as compared to the total degeneracy $N$ of the level. Therefore, at $t \sim N^{3 / 2}\left(8\left\langle V^{2}\right\rangle\right)^{-1 / 2}$, the system no longer conforms to the model employed. At this time, the population recurrences to energy eigenstates become important. Later on, we will consider this situation in more detail.

### 8.1.2 Non-Degenerate Bands With Random Coupling

We now consider the case of non-degenerate bands, as shown in Fig. 8.1(b). The three-term recurrence relation (8.4) instead of the operators $X(\varepsilon)$ and $X(\xi)$ of (8.6) now contains the operators (4.33) for non-degenerate bands, which are slightly modified by the presence of two neighbors of each band and read

$$
\begin{align*}
X_{n, k}(\varepsilon) & =\frac{1}{\varepsilon-\Delta_{n, k}+i \pi\left\langle V_{n+1}^{2}\right\rangle g_{n+1}+i \pi\left\langle V_{n}^{2}\right\rangle g_{n-1}} \\
X_{n, k}(\xi) & =\frac{1}{\xi-\Delta_{n, k}-i \pi\left\langle V_{n+1}^{2}\right\rangle g_{n+1}-i \pi\left\langle V_{n}^{2}\right\rangle g_{n-1}} \tag{8.20}
\end{align*}
$$

in the general case of bands of different density of states $g_{n}$ and different mean square interactions $\left\langle V_{n}^{2}\right\rangle$ coupling the band $n$ with the band $n-1$. Substitution of (8.20) into (8.4), after summation over $k$ yields

$$
\begin{align*}
\rho_{n}(\varepsilon, \xi)= & \frac{2 \pi i \rho_{n-1}(\varepsilon, \xi)\left\langle V_{n}^{2}\right\rangle g_{n}}{\varepsilon-\xi+2 i \pi\left\langle V_{n+1}^{2}\right\rangle g_{n+1}+2 i \pi\left\langle V_{n}^{2}\right\rangle g_{n-1}} \\
& +\frac{2 \pi i \rho_{n+1}(\varepsilon, \xi)\left\langle V_{n+1}^{2}\right\rangle g_{n}}{\varepsilon-\xi+2 i \pi\left\langle V_{n+1}^{2}\right\rangle g_{n+1}+2 i \pi\left\langle V_{n}^{2}\right\rangle g_{n-1}}+\delta_{n}^{0} X_{0,0}(\varepsilon) X_{0,0}(\xi) . \tag{8.21}
\end{align*}
$$

Introducing the variables $\zeta=\varepsilon-\xi, \eta=(\varepsilon+\xi) / 2$ and integrating this equation over $d \eta$ one finds

$$
\begin{align*}
\rho_{n}(\zeta)= & \frac{2 i \pi \rho_{n-1}(\zeta)\left\langle V_{n}^{2}\right\rangle g_{n}}{\zeta+2 i \pi\left\langle V_{n+1}^{2}\right\rangle g_{n+1}+2 i \pi\left\langle V_{n}^{2}\right\rangle g_{n-1}} \\
& +\frac{2 i \pi \rho_{n+1}(\zeta)\left\langle V_{n+1}^{2}\right\rangle g_{n}}{\zeta+2 i \pi\left\langle V_{n+1}^{2}\right\rangle g_{n+1}+2 i \pi\left\langle V_{n}^{2}\right\rangle g_{n-1}}+\frac{2 i \pi \delta_{n}^{0}}{\zeta+2 i \pi\left\langle V_{1}^{2}\right\rangle g_{1}+2 i \pi\left\langle V_{0}^{2}\right\rangle g_{-1}} . \tag{8.22}
\end{align*}
$$

for $\rho_{n}(\zeta)=\int \rho_{n}(\zeta, \eta) d \eta / 2 \pi$. This immediately yields

$$
\begin{align*}
\zeta \rho_{n}(\zeta)= & 2 i \pi\left\langle V_{n}^{2}\right\rangle\left[\rho_{n-1}(\zeta) g_{n}-\rho_{n}(\zeta) g_{n-1}\right] \\
& +2 i \pi\left\langle V_{n+1}^{2}\right\rangle\left[\rho_{n+1}(\zeta) g_{n}-\rho_{n}(\zeta) g_{n+1}\right]+\delta_{n}^{0} \delta_{k}^{0} \tag{8.23}
\end{align*}
$$

and after inverse Fourier transformation over $\zeta$ results in the master equation

$$
\begin{align*}
\dot{\rho}_{n}(t)= & 2 \pi\left\langle V_{n}^{2}\right\rangle\left[\rho_{n-1}(t) g_{n}-\rho_{n}(t) g_{n-1}\right] \\
& +2 \pi\left\langle V_{n+1}^{2}\right\rangle\left[\rho_{n+1}(t) g_{n}-\rho_{n}(t) g_{n+1}\right]+\delta_{n}^{0} \delta(t) \tag{8.24}
\end{align*}
$$

where the last term containing the Dirac $\delta$-function allows for the initial condition $\rho_{n}(t=0)=\delta_{n}^{0}$. One can interpret this term as a unit population instantaneously injected into the band $n=0$ at time $t=0$.

For the population per quantum state $\widetilde{\rho}_{n}(t)=\rho_{n}(t) / g_{n},(8.24)$ takes a more symmetric form

$$
\begin{equation*}
\dot{\tilde{\rho}}_{n}(t)=D_{n}^{-}\left[\widetilde{\rho}_{n-1}(t)-\widetilde{\rho}_{n}(t)\right]+D_{n}^{+}\left[\widetilde{\rho}_{n+1}(t)-\widetilde{\rho}_{n}(t)\right]+\frac{\delta_{n}^{0} \delta(t)}{g_{0}} \tag{8.25}
\end{equation*}
$$

where $D_{n}^{+}=2 \pi g_{n+1}\left\langle V_{n+1}^{2}\right\rangle$ and $D_{n}^{-}=2 \pi g_{n-1}\left\langle V_{n}^{2}\right\rangle$ denote the transition probability to the neighboring bands with the larger and smaller $n$, respectively. An attempt to solve this equation analytically leads to the same problems that have already been discussed in the context of the Schrödinger equation (7.1) for the relay of isolated levels - there are only a small number of
particular cases of dependences $D_{n}^{-}$and $D_{n}^{+}$where the problem possesses an exact solution. However, since the eigenfunctions of (8.25) decay exponentially in time, the approximate WKB calculations describing the time evolution of the "wavepackets" turn out to be rather accurate, and applicable to most cases.

## Constant Transition Rates

For the case of $n$-independent coefficients $D_{n}^{-}$and $D_{n}^{+}$, solution of (8.25) can be found with the help the generating function

$$
\begin{equation*}
\Phi(x, t)=\sum_{n=-\infty}^{\infty} e^{i n x} \widetilde{\rho}_{n}(t) \tag{8.26}
\end{equation*}
$$

One finds the differential equation

$$
\begin{equation*}
\dot{\Phi}(x, t)=\left[D^{-}\left(e^{i x}-1\right)+D^{+}\left(e^{-i x}-1\right)\right] \Phi(x, t)+\frac{\delta(t)}{g_{0}} \tag{8.27}
\end{equation*}
$$

for this function which immediately yields

$$
\begin{equation*}
\Phi(x, t)=\frac{1}{g_{0}} \exp \left[t D^{-}\left(e^{i x}-1\right)+t D^{+}\left(e^{-i x}-1\right)\right] \tag{8.28}
\end{equation*}
$$

From the generating function (8.28) one finds

$$
\begin{align*}
\widetilde{\rho}_{n}(t) & =\frac{1}{2 \pi} \int_{-\pi}^{\pi} d x e^{-i n x} \Phi(x, t)=\frac{\left.e^{-t\left(D^{-}+D^{+}\right.}\right)}{2 \pi i g_{0}} \int_{C} \frac{d y}{y^{n+1}} \exp \left[t D^{-} y+t D^{+} \frac{1}{y}\right] \\
& \left.=e^{-t\left(D^{-}+D^{+}\right.}\right) \frac{I_{n}\left(2 t \sqrt{D^{-} D^{+}}\right)}{g_{0}\left(D^{+} / D^{-}\right)^{n / 2}} \tag{8.29}
\end{align*}
$$

where the replacement $e^{i x} \rightarrow y$ has been done. The contour $C$ is a loop around $y=0$, and $I_{n}(x)$ is a modified Bessel function. Note that the requirements $D_{n}^{-}=$const and $D_{n}^{+}=$const imply that $g_{n}=$ const and $\left\langle V_{n}^{2}\right\rangle=$ const, and consequently one finds $D_{n}^{-}=D_{n}^{+}=D$, which yields for the total populations

$$
\begin{equation*}
\rho_{n}(t)=e^{-2 D t} I_{n}(2 D t) \tag{8.30}
\end{equation*}
$$

Let us find the asymptotic form of this equation for large $t$ when the main contribution to the integral (8.29) comes from the vicinity of the saddle point $y=1$, that is $x=0$. By developing the argument of the exponent in (8.28) over $x$, up to second order and by extending the integration limits for $x$ to $\pm \infty$ one finds the distribution

$$
\begin{align*}
\rho_{n}(t) & =e^{-2 D t} I_{n}(2 D t) \\
& \simeq \frac{1}{2 \pi} \int_{-\infty}^{\infty} d x \exp \left[-t D x^{2}-i n x\right]=\frac{1}{\sqrt{2 \pi D t}} \exp \left[-\frac{n^{2}}{4 D t}\right] \tag{8.31}
\end{align*}
$$

coinciding with the classical diffusion expression. We note that at $t>1 / D$ for the population distribution $\rho_{n}(\Delta, t)$ over the band one finds the Lorentzian profile (4.39) with $W=2 D$ by considering the poles of $\rho_{n, k}(\zeta, \eta)$ on the real axis of $\zeta$ by analogy to (4.38), which yields

$$
\begin{equation*}
\rho_{n}(\Delta, t) \simeq \rho_{n}(t) \frac{2 D}{\pi\left(4 D^{2}+\Delta^{2}\right)} \tag{8.32}
\end{equation*}
$$

In the context of (8.32), it is expedient to discuss the applicability of the approach based on the master equation (8.25) for asymptotically long times. As for the case of degenerate levels, the validity of the approach is limited in time by the assumption of the tree-like structure of diagrams, which implies that self-intersections of the diagrams are unlikely. This assumption fails when the typical number of returns of the random trajectory $\sqrt{D t}$ becomes of the order of the total number $g D$ of band levels in the populated domain of width $D$, suggested by the distribution (8.31). This occurs for $t \sim g^{2} D$, when the distribution over the band has a width $g D$ coinciding with the number of populated levels in each band.

## Linearly Increasing Transition Rates

Equation (8.25) can also be solved exactly for linearly changing coefficients $D_{n}^{+}=2 \pi g_{n+1}\left\langle V_{n+1}^{2}\right\rangle$ and $D_{n}^{-}=2 \pi g_{n-1}\left\langle V_{n}^{2}\right\rangle$. This variation can result from a linear dependence of the state density $g_{n}=g(1+\alpha n)$ on the band number $n$ or from the dependence of the mean square coupling $\left\langle V_{n}^{2}\right\rangle=V^{2}(1+\alpha n)$ on $n$. A slight difference between these two cases manifests itself in a difference between the coefficients $D_{n}^{-}$: the coefficient for the first case $D_{n+1}^{-}=2 \pi g\left\langle V^{2}\right\rangle$ $(1+\alpha n)$ coincides with the coefficient $D_{n}^{-}$for the second case, whereas the coefficients $D_{n}^{+}$are identical. Equation (8.25) takes the form

$$
\begin{align*}
\dot{\tilde{\rho}}_{n}(\tau)= & \left(1-\alpha^{\prime}+\alpha n\right)\left[\widetilde{\rho}_{n-1}(\tau)-\widetilde{\rho}_{n}(\tau)\right] \\
& +(1+\alpha+\alpha n)\left[\widetilde{\rho}_{n+1}(\tau)-\widetilde{\rho}_{n}(\tau)\right]+\frac{\delta_{n}^{0} \delta(\tau)}{g} \tag{8.33}
\end{align*}
$$

where $\tau=2 \pi g\left\langle V^{2}\right\rangle t$, and either $\alpha^{\prime}=\alpha$ for the first case, or $\alpha^{\prime}=0$ for the second case.

For the generating function (8.26) one finds the equation

$$
\begin{align*}
\dot{\Phi}(x, \tau)= & {\left[\left(1-\alpha^{\prime}\right)\left(e^{i x}-1\right)+(1+\alpha)\left(e^{-i x}-1\right)\right] \Phi(x, \tau) } \\
& -i \alpha \frac{\partial}{\partial x}\left\{\left[\left(e^{i x}-1\right)+\left(e^{-i x}-1\right)\right] \Phi(x, \tau)\right\}+\frac{\delta(\tau)}{g} . \tag{8.34}
\end{align*}
$$

where linear dependence on $n$ corresponds to the derivative $-i \partial / \partial x$. Solution of this first-order linear differential equation in two dimensions can be found by the standard method of characteristics. This reads

$$
\begin{equation*}
\Phi(x, \tau)=\frac{1}{g} e^{-i x / \alpha}\left(e^{i x}+\alpha \tau-e^{i x} \alpha \tau\right)^{1 / \alpha}\left(1+\alpha \tau-e^{i x} \alpha \tau\right)^{-1+\alpha^{\prime} / \alpha-1 / \alpha} \tag{8.35}
\end{equation*}
$$

and yields the population
$\widetilde{\rho}_{n}(\tau)=\frac{1}{2 \pi g} \int_{0}^{2 \pi} d x e^{-i n x-i x / \alpha}\left(e^{i x}+\alpha \tau-e^{i x} \alpha \tau\right)^{1 / \alpha}\left(1+\alpha \tau-e^{i x} \alpha \tau\right)^{-1+\alpha^{\prime} / \alpha-1 / \alpha}$.
By making the replacement $e^{i x} \rightarrow y$ and by deformation of the integration contour one can write this integral in the form

$$
\begin{equation*}
\widetilde{\rho}_{n}(\tau)=\frac{-i}{2 \pi g} \int_{C_{0}} d y y^{-n-1 / \alpha}(y+\alpha \tau-y \alpha \tau)^{1 / \alpha}(1+\alpha \tau-y \alpha \tau)^{-1+\alpha^{\prime} / \alpha-1 / \alpha} \tag{8.37}
\end{equation*}
$$

where the integration contour $C_{0}$ circumvents the zeros of two last factors of the integrand. This three-center integral resembles the integral representation of the hypergeometric function ${ }_{2} F_{1}(a, b ; c ; X)$, and indeed, after the replacement $y \rightarrow\left(\alpha^{2} t^{2}-z\right) /\left(\alpha^{2} t^{2}-\alpha t\right)$ the integral over $d z$ takes the form of this integral representation and yields

$$
\begin{equation*}
\widetilde{\rho}_{n}(\tau)=\frac{\Gamma(n-1 / \alpha)(-\alpha \tau)^{n}{ }_{2} F_{1}\left(n-\left(1-\alpha^{\prime}\right) / \alpha, n-1 / \alpha ; n+1 ; \alpha^{2} \tau^{2}\right)}{g(1+\alpha \tau)^{1+2 / \alpha-\alpha^{\prime} / \alpha-n} \Gamma(-1 / \alpha) n!} \tag{8.38}
\end{equation*}
$$

which for large $n$ and $\tau$ has the asymptotic form

$$
\begin{equation*}
\widetilde{\rho}_{n}(\tau) \simeq \frac{1}{2 \pi g} \int_{-\infty}^{\infty} d x e^{-i n x-i x / \alpha}(1-i x \alpha \tau)^{-1} \simeq \frac{e^{-n / \alpha \tau}}{\alpha \tau g} \tag{8.39}
\end{equation*}
$$

in the case $\alpha^{\prime}=\alpha$. For $\alpha^{\prime}=0$ the asymptotic behavior differs by the preexponential factor.

In Fig. 8.6 we illustrate the dependence of the population per state on the band number. One sees that the population distribution is limited in the domain of $n$, where the coefficients $D_{n}^{+}$and $D_{n}^{-}$are positive. For the case of linearly dependent coupling, the distribution tends to zero when $n$ approaches the band for which one of the coefficients $D$ equals zero. For the linearly dependent state density at this point, the distribution function is discontinuous. Both functions manifest asymptotic behavior (8.39).

## WKB Approximation

For smooth dependencies of the kinetic coefficients $D_{n}^{+}$and $D_{n}^{-}$on the band number, one can employ the WKB approximation by analogy to that considered on p. 326. With the help of the finite shift operators (5.15) one can write the continuous analog of (8.25) in the form


Fig. 8.6. Dependence of the population per state on the band number for the linearly changing transition probability. (a) Dependence of the population on the band number $n$ and on the scaled time $\tau=2 \pi g\left\langle V^{2}\right\rangle t$ in the case of constant state density and linearly increasing coupling ( $\alpha=0.2, \alpha^{\prime}=0$ ). (b) Population distributions for the linearly increasing state density $g_{n}=g(1+0.2 n)$ and for linearly increasing mean square coupling $\left\langle V_{n}^{2}\right\rangle=\left\langle V^{2}\right\rangle(1+0.2 n)$ are shown by dotted lines 1 and 2, respectively.

$$
\begin{equation*}
\dot{\tilde{\rho}}(n, t)=\left\{D^{-}(n)\left[\exp \left(-\frac{\partial}{\partial n}\right)-1\right]+D^{+}(n)\left[\exp \left(\frac{\partial}{\partial n}\right)-1\right]\right\} \widetilde{\rho}(n, t), \tag{8.40}
\end{equation*}
$$

which requires separate consideration of the initial condition. Smooth dependence of the positive coefficients results in a smooth distribution function, for which one can expand exponents into a Taylor series up to second order and obtain

$$
\begin{equation*}
\dot{\tilde{\rho}}(n, t)=\left[D^{+}(n)-D^{-}(n)\right] \frac{\partial \widetilde{\rho}(n, t)}{\partial n}+\left[D^{+}(n)+D^{-}(n)\right] \frac{\partial^{2} \widetilde{\rho}(n, t)}{2 \partial n^{2}} \tag{8.41}
\end{equation*}
$$

The WKB ansatz $\widetilde{\rho}(n, t)=\int d \lambda g(\lambda) \exp [-\lambda t+S(n, \lambda)]$ yields in the same order of approximation

$$
\begin{equation*}
-\lambda=\Pi(n) \frac{\partial S(n, \lambda)}{\partial n}+D(n)\left[\frac{\partial S(n, \lambda)}{\partial n}\right]^{2} \tag{8.42}
\end{equation*}
$$

where

$$
\begin{align*}
& \Pi(n)=\left[D^{+}(n)-D^{-}(n)\right] \\
& D(n)=\frac{1}{2}\left[D^{+}(n)+D^{-}(n)\right], \tag{8.43}
\end{align*}
$$

whereas the function $g(\lambda)$ has to be found from the initial condition at $t=0$.
Solution of the quadratic equation (8.42) allows one to find the exponential factor

$$
\begin{equation*}
S(n, \lambda)=\int_{0}^{n} d n^{\prime}\left[\frac{-\Pi\left(n^{\prime}\right)}{2 D\left(n^{\prime}\right)} \pm \frac{\sqrt{\Pi^{2}\left(n^{\prime}\right)+4 \lambda D\left(n^{\prime}\right)}}{2 D\left(n^{\prime}\right)}\right] \tag{8.44}
\end{equation*}
$$

where the positive sign, which gives a physically meaningful solution, has to be chosen. For long times only small $\lambda$ are important, and therefore one can make use of the Taylor expansion of the right-hand side of (8.44), that is

$$
\begin{equation*}
S(n, \lambda)=\int_{0}^{n} d n^{\prime}\left[\frac{\lambda}{\Pi\left(n^{\prime}\right)}-\frac{\lambda^{2} D\left(n^{\prime}\right)}{\Pi^{3}\left(n^{\prime}\right)}\right] \tag{8.45}
\end{equation*}
$$

Therefore the WKB ansatz takes the form

$$
\begin{equation*}
\widetilde{\rho}(n, t) \simeq \int d \lambda g(0) \exp \left[-\lambda\left(t-\int_{0}^{n} \frac{d n^{\prime}}{\Pi\left(n^{\prime}\right)}\right)-\lambda^{2} \int_{0}^{n} d n^{\prime} \frac{D\left(n^{\prime}\right)}{\Pi^{3}\left(n^{\prime}\right)}\right] . \tag{8.46}
\end{equation*}
$$

Integration over $d \lambda$ yields

$$
\begin{equation*}
\widetilde{\rho}(n, t) \simeq \sqrt{2 \pi \int_{0}^{n} \frac{D\left(n^{\prime}\right) d n^{\prime}}{\Pi^{3}\left(n^{\prime}\right)}} g(0) \exp \left[-\frac{\left(t-\int_{0}^{n} \frac{d n^{\prime}}{\Pi\left(n^{\prime}\right)}\right)^{2}}{4 \int_{0}^{n} d n^{\prime} \frac{D\left(n^{\prime}\right)}{\Pi^{3}\left(n^{\prime}\right)}}\right] . \tag{8.47}
\end{equation*}
$$

If the dependencies $\Pi(n)$ and $D\left(n^{\prime}\right)$ are such that for large $n$

$$
\begin{equation*}
\int_{0}^{n} \frac{d n^{\prime}}{\Pi\left(n^{\prime}\right)} \gg\left(\int_{0}^{n} d n^{\prime} \frac{D\left(n^{\prime}\right)}{\Pi^{3}\left(n^{\prime}\right)}\right)^{1 / 2} \tag{8.48}
\end{equation*}
$$

then at a given time moment $t$, expression (8.47) corresponds to a distribution, localized around the band number $n_{0}$ for which

$$
\begin{equation*}
\int_{0}^{n_{0}} \frac{d n^{\prime}}{\Pi\left(n^{\prime}\right)}=t \tag{8.49}
\end{equation*}
$$

while the distribution width is

$$
\begin{equation*}
\Delta n \sim \Pi\left(n_{0}\right)\left(\int_{0}^{n_{0}} d n^{\prime} \frac{D\left(n^{\prime}\right)}{\Pi^{3}\left(n^{\prime}\right)}\right)^{1 / 2} \tag{8.50}
\end{equation*}
$$

Note that if the relation (8.48) is not valid, the width of the distribution $\Delta n$ exceeds the displacement of the center of the distribution $n_{0}$ and the truncation of the Taylor expansion (8.45) at the second-order terms might not be justified. In this case (8.47) might not be very accurate. In the latter case one has to perform an exact integration over $d n^{\prime} d \lambda$ employing (8.44).

### 8.1.3 Correlated Coupling

Now we consider the case of (8.1) when the coupling matrices $\widehat{V}_{n}^{n+1}$ with the matrix elements $V_{n, k}^{n+1, m}$ connecting each level of a band with the levels of the neighboring band have low ranks $R \ll N$. Here $N$ denotes a typical number of levels in each band. In other words, the matrix elements $\widehat{V}_{n}^{n+1}=\sum_{j=1}^{R}\left|C_{n}^{(j)}\right\rangle \mathcal{V}^{(j)}\left\langle C_{n+1}^{(j)}\right|$ are strongly correlated, since each of the matrices is a sum of tensor products for a small number $R$ of eigenvectors. Earlier, on p. 160 we considered such a situation for two-band systems. Now, considering multiband problems, we concentrate on the extreme, when the rank of the matrix is one, and therefore the matrix elements have a factorized form $V_{n, k}^{n+1, m}=C_{n, k}^{\prime} \mathcal{V}_{n} C_{n+1, m}$ determined by the single non-zero eigenvalue $\mathcal{V}_{n}$ and the corresponding left $C_{n, k}^{\prime}$ and right $C_{n+1, m}$ eigenvectors normalized to unity, for each band number $n$. We assume that the components $C_{n, k}^{\prime}$ and $C_{n+1, m}$ of the vectors are real, although a straightforward generalization of the complex values of these coefficients can also be performed. We note that the eigenvalue $\mathcal{V}_{n}$ may be physically interpreted as the cooperative matrix element of the transition between band $n$ and band $n+1 . A$ typical matrix element $V$, coupling two levels of neighboring bands, is of the order of $\mathcal{V} / N$.

Substitution of this particular form of $V_{n, k}^{n+1, m}$ into (8.1) after Fourier transformation results in

$$
\begin{align*}
\left(\varepsilon-E_{n, k}\right) \psi_{n, k}(\varepsilon)= & i \psi_{n, k}(t=0)+C_{n, k}^{\prime} \mathcal{V}_{n} \sum_{m} C_{n+1, m} \psi_{n+1, m}(\varepsilon) \\
& +C_{n, k} \mathcal{V}_{n-1} \sum_{m} C_{n-1, m}^{\prime} \psi_{n-1, m}(\varepsilon), \tag{8.51}
\end{align*}
$$

where the term $i \psi_{n, k}(t=0)$ on the right-hand side correspond $s$ to the initial condition. For the combinations

$$
\begin{align*}
& \Phi_{n}(\varepsilon)=\sum_{k} C_{n, k} \psi_{n, k}(\varepsilon) \\
& \bar{\Phi}_{n}(\varepsilon)=\sum_{k} C_{n, k}^{\prime} \psi_{n, k}(\varepsilon) \tag{8.52}
\end{align*}
$$

(8.51) yields

$$
\begin{align*}
& \Phi_{n}=\mathcal{V}_{n} \Phi_{n+1} \sum_{k} \frac{C_{n, k} C_{n, k}^{\prime}}{\varepsilon-E_{n, k}}+\mathcal{V}_{n-1} \bar{\Phi}_{n-1} \sum_{k} \frac{C_{n, k} C_{n, k}}{\varepsilon-E_{n, k}}+i U_{n} \\
& \bar{\Phi}_{n}=\mathcal{V}_{n} \Phi_{n+1} \sum_{k} \frac{C_{n, k}^{\prime} C_{n, k}^{\prime}}{\varepsilon-E_{n, k}}+\mathcal{V}_{n-1} \bar{\Phi}_{n-1} \sum_{k} \frac{C_{n, k}^{\prime} C_{n, k}}{\varepsilon-E_{n, k}}+i \bar{U}_{n} \tag{8.53}
\end{align*}
$$

where

$$
\begin{align*}
& U_{n}(\varepsilon)=\sum_{k} \frac{C_{n, k} \psi_{n, k}(t=0)}{\varepsilon-E_{n, k}} \\
& \bar{U}_{n}(\varepsilon)=\sum_{k} \frac{C_{n, k}^{\prime} \psi_{n, k}(t=0)}{\varepsilon-E_{n, k}} \tag{8.54}
\end{align*}
$$

In Figs. 8.1(c) and (d) one sees two different cases of the system with factorized interaction. The case of Fig. 8.1(c) corresponds to the situation where the left $\left|C_{n}^{(j)}\right\rangle$ and the conjugate of the right eigenstate $\left|C_{n+1}^{(j)}\right\rangle$ of the coupling matrices coincide, and therefore the interaction couples these "bright" states of neighboring bands directly. In this case the population propagates along this relay of isolated "bright" levels, decaying to the other "dark" states of the bands, by analogy to the situation discussed on p. 170 for the two-band system. In the case of Fig. 8.1(d) the left and right eigenvectors for the coupling are different. This means that the state which is "bright" for the transition $n-1 \rightarrow n$ does not coincide with the "bright" state for the transition $n \rightarrow n+1$. The matrix element of the resolvent coupling these two states determines the efficiency of the transition, by analogy to the quantity $W_{1}^{0}$ in (4.114) introduced for the problem of two states interacting via the band shown in Fig. 4.22.

For statistically identical bands, and for times shorter as compared to the typical density of states of one band, the sums

$$
\begin{array}{ll}
W_{0}^{1}=\mathcal{V}_{n} \sum_{k} \frac{C_{n, k} C_{n, k}^{\prime}}{\varepsilon-E_{n, k}}, & W_{0}^{0}=\mathcal{V}_{n-1} \sum_{k} \frac{C_{n, k} C_{n, k}}{\varepsilon-E_{n, k}} \\
W_{1}^{1}=\mathcal{V}_{n} \sum_{k} \frac{C_{n, k}^{\prime} C_{n, k}^{\prime}}{\varepsilon-E_{n, k}}, & W_{1}^{0}=\mathcal{V}_{n-1} \sum_{k} \frac{C_{n, k}^{\prime} C_{n, k}}{\varepsilon-E_{n, k}}, \tag{8.55}
\end{array}
$$

do not depend on the index $n, W_{0}^{1}=W_{1}^{0}=W$ and therefore (8.53) takes the form

$$
\begin{align*}
& \Phi_{n}=\Phi_{n+1} W+\bar{\Phi}_{n-1} W_{0}^{0}+i U_{n} \\
& \bar{\Phi}_{n}=\Phi_{n+1} W_{1}^{1}+\bar{\Phi}_{n-1} W+i \bar{U}_{n} \tag{8.56}
\end{align*}
$$

and results in the equation

$$
\begin{align*}
& \Phi(x)=e^{-i x} \Phi(x) W+e^{i x} \bar{\Phi}(x) W_{0}^{0}+i U(x) \\
& \bar{\Phi}(x)=e^{-i x} \Phi(x) W_{1}^{1}+e^{i x} \bar{\Phi}(x) W+i \bar{U}(x) \tag{8.57}
\end{align*}
$$

for the generating functions

$$
\begin{array}{ll}
\Phi(x)=\sum_{n} e^{i n x} \Phi_{n}, & \bar{\Phi}(x)=\sum_{n} e^{i n x} \bar{\Phi}_{n} \\
U(x)=\sum_{n} e^{i n x} U_{n}, & \bar{U}(x)=\sum_{n} e^{i n x} \bar{U}_{n} \tag{8.58}
\end{array}
$$

Solution of (8.58) yields for the functions (8.52)

$$
\begin{align*}
& \Phi_{n}=\frac{i}{2 \pi} \int_{-\pi}^{\pi} \frac{\left(1-e^{i x} W\right) U(x)-e^{-i x} W_{1}^{1} \bar{U}(x)}{1-e^{i x} W-e^{-i x} W} e^{-i n x} d x \\
& \bar{\Phi}_{n}=\frac{i}{2 \pi} \int_{-\pi}^{\pi} \frac{\left(1-e^{-i x} W\right) \bar{U}(x)-e^{i x} W_{0}^{0} U(x)}{1-e^{i x} W-e^{-i x} W} e^{-i n x} d x \tag{8.59}
\end{align*}
$$

where the argument $\varepsilon$ is implicit.

If all of the population is initially in the band $n=0$, then the generating function of the initial conditions $U(x)$ does not depend on the variable $x$. In this case the integrands (8.59) only have singularities at the points where $\cos x=1 / 2 W$ and the integration yields

$$
\begin{align*}
& \Phi_{n}=U \delta_{n}^{0} \pm \frac{U W-W_{1}^{1} \bar{U}}{\sqrt{1-4 W^{2}}} e^{\mp i(n+1) \arccos (1 / 2 W)} \\
& \bar{\Phi}_{n}=\bar{U} \delta_{n}^{0} \pm \frac{\bar{U} W-W_{0}^{0} U}{\sqrt{1-4 W^{2}}} e^{\mp i(n-1) \arccos (1 / 2 W)} \tag{8.60}
\end{align*}
$$

where the choice of the sign + or - depends on the sign of the imaginary part of $\arccos 1 / 2 W$ and the sign of the number $n$ in such a way, that the result converges for $|n| \rightarrow \infty$. Substitution of (8.60) into (8.51) gives the probability amplitude

$$
\begin{align*}
\psi_{n, k}(\varepsilon)= & \frac{C_{n, k}^{\prime} \mathcal{V}\left(U W-W_{1}^{1} \bar{U}\right)}{\left(\varepsilon-E_{n, k}\right) \sqrt{1-4 W^{2}}} e^{\mp i(n+2) \arccos (1 / 2 W)} \\
& +\frac{C_{n, k} \mathcal{V}\left(\bar{U} W-W_{0}^{0} U\right)}{\left(\varepsilon-E_{n, k}\right) \sqrt{1-4 W^{2}}} e^{\mp i(n-2) \arccos (1 / 2 W)} . \tag{8.61}
\end{align*}
$$

for the initially non-populated state.
For further calculations one needs to make a particular choice of the real coefficients $C_{n, k}$ and $C_{n, k}^{\prime}$. We take the simplest model and assume the distribution

$$
\begin{align*}
\left(C_{n, k}\right)^{2} & =\left(C_{n, k}^{\prime}\right)^{2}=\frac{N}{\pi\left(k^{2}+N^{2}\right)} \\
C_{n, k} C_{n, k}^{\prime} & =\frac{\mu N}{\pi\left(k^{2}+N^{2}\right)} \tag{8.62}
\end{align*}
$$

for the squares and the product of these coefficients by analogy to (4.138). Here the factor $\mu$ entering the last equation allows for the fact that the left and the right eigenvectors of the coupling matrix can be different, although they have an identical distribution for the absolute values. The difference manifests itself in the statistics of the coefficients, and we assume that the product $C_{n, k}^{\prime} C_{n, k}$ averaged over a narrow strip of energy, near the energy position $E_{n, k}$ differs only by an energy-independent factor $\mu$ from the product $C_{n, k} C_{n, k}$ averaged over the same strip. We have already encountered such a situation in (4.114), considering the transition between two isolated levels which occurs via a band. Equation (8.62) also implies that the bands are infinitely broad, although the levels located a large distance from the energy position $E_{n, k} \simeq 0$ are weakly coupled to the levels of neighboring bands, since they are weakly represented in the eigenvectors $\left|C_{n}\right\rangle$ of the coupling matrices.

Substitution of (8.62) into (8.55) yields

$$
\begin{align*}
W_{0}^{0} & =W_{1}^{1}=\mathcal{V} \sum_{k} \frac{1}{\varepsilon-k \delta} \frac{N}{\pi\left(k^{2}+N^{2}\right)}=\frac{\mathcal{V}}{\varepsilon+i N \delta} \\
W & =\mathcal{V} \sum_{k} \frac{1}{\varepsilon-k \delta} \frac{\mu N}{\pi\left(k^{2}+N^{2}\right)}=\frac{\mathcal{V} \mu}{\varepsilon+i N \delta} \tag{8.63}
\end{align*}
$$

for $\mathcal{V}_{n}$ independent of $n$, and for $E_{n, k}=\delta k$, where $\delta$ is a typical energy spacing among neighboring levels in the same band. As we have already seen in Chap. 3, the detailed structure of the spectrum is not important for time intervals $t \ll 1 / \delta$ shorter then the recurrence period, and this spectrum can be considered as an equidistant one, whereas the sums can be replaced by integrals.

The main effects associated with correlated matrix elements are seen already for the simplest case of the initial condition

$$
\begin{equation*}
\psi_{n, k}(t=0)=C_{n, k}, \tag{8.64}
\end{equation*}
$$

which means that the state $\left|C_{n}\right\rangle$ corresponding to the eigenvector of the coupling matrix is populated at $t=0$. Therefore comparing (8.54) with allowance of (8.64) and (8.55) we find

$$
\begin{align*}
& U(\varepsilon)=\sum_{k} \frac{C_{n, k} C_{n, k}}{\varepsilon-E_{n, k}}=\frac{W}{\mathcal{V}} \\
& \bar{U}(\varepsilon)=\sum_{k} \frac{C_{n, k}^{\prime} C_{n, k}}{\varepsilon-E_{n, k}}=\frac{\mu W}{\mathcal{V}} \tag{8.65}
\end{align*}
$$

and substitution of Eqs. $(8.63,8.65)$ into (8.61) yields

$$
\begin{align*}
\psi_{n, k}(\varepsilon)= & \frac{\left(1-\mu^{2}\right)(\mathcal{V} \mu)^{2}}{(i N-k) \delta} \sqrt{\frac{N}{\pi\left(k^{2}+N^{2}\right)\left[(\varepsilon+i N \delta)^{2}-4(\mathcal{V} \mu)^{2}\right]}} \\
& \left(\frac{1}{\varepsilon-k \delta}-\frac{1}{\varepsilon-i N \delta}\right) \exp \left[\mp i(n-2) \arccos \left(\frac{\varepsilon+i N \delta}{2 \mathcal{V} \mu}\right)\right] \tag{8.66}
\end{align*}
$$

This expression, as a function of $\varepsilon$, has four singularities: two poles at the points $\varepsilon=k \delta$ and $\varepsilon=-i N \delta$, and two branching points at $\varepsilon=-i N \delta \pm 2 \mathcal{V} \mu$. Only the first point is located on the real axis and therefore is important at times $t>1 / N \delta$, whereas the other points yield exponentially decreasing contributions responsible for transitional dynamics resembling that of (7.97) for the relay of isolated levels. Contribution of the pole at the real axis yields the stationary distribution

$$
\begin{align*}
\psi_{n, k}(t \rightarrow \infty) & =\frac{\pi i\left(\mu^{2}-1\right)}{2 \mathcal{G}(i N+k)} \\
& \sqrt{\frac{N}{\pi\left(k^{2}+N^{2}\right)}} \frac{\exp [\mp i(n-2) \arccos (\mathcal{G}(k+i N))]}{\sqrt{\mathcal{G}^{2}(k+i N)^{2}-1}} \tag{8.67}
\end{align*}
$$

which depends only on one dimensionless parameter $\mathcal{G}=\delta / 2 \mathcal{V} \mu$ and results in the expression

$$
\begin{align*}
\rho_{n, k}= & \left(\frac{\left(1-\mu^{2}\right)}{2 \mathcal{G}}\right)^{2} \frac{N \pi}{\left(N^{2}+k^{2}\right)^{2}} \\
& \frac{\exp [-2(n-2)|\operatorname{Re} \arccos (\mathcal{G}(k+i N))|]}{\sqrt{\left[\mathcal{G}^{2}(k-i N)^{2}-1\right]\left[\mathcal{G}^{2}(k+i N)^{2}-1\right]}} \tag{8.68}
\end{align*}
$$

for the population distribution.



Fig. 8.7. Population distribution (not to scale) over the bands for the case of correlated coupling of band levels. The typical number of interacting levels in each band is $N=1000$. Exponential factors have been introduced in order to illustrate the energy distribution of the population for large $n$. (a) For strong interaction $\mathcal{G} N \ll 1$ one sees the population peaks detuned from the center, resembling that shown in Fig. 4.25. (b) Weak interaction. Population distribution over the energy remains Lorentzian, whereas the total population distribution amongst the bands is localized practically only on one, the initially populated band.

From (8.68), it is clear that at $k=0$ in the center of the distribution over the bands, the population as a function of $n$ manifests an exponential behavior that is $\rho_{n, 0} \sim \exp [-2 n|\operatorname{Re} \beta|]$, where $\beta=\arccos (i \mathcal{G} N)$. The parameter $\mathcal{G} N=N \delta / 2 \mathcal{V} \mu$ may be physically interpreted as the ratio of the effective band width $N \delta$ and the cooperative matrix element of the interaction $\mathcal{V}$. We have already encountered a similar parameter $1 / \sqrt{w}=\sqrt{N \delta_{0} M \delta_{1}} / \mathcal{V}$ that is responsible for the population distribution width in the two-band system shown in Fig. 4.25. When this parameter is large one finds $|\operatorname{Re} \beta| \simeq \log (\mathcal{G} N)$ and the population distribution over the bands behaves as $\rho_{n, 0} \sim(\mathcal{G} N)^{-2 n}$, decreasing rapidly with the increase of $n$. As we have seen on p. 387, a typical matrix element $V$ which couples two levels of neighboring bands is of order
$\mathcal{V} / N$, and therefore the parameter $\mathcal{G} N$ is of the order of the ratio of the typical detuning $\delta / 2$ between the two closest levels of the neighboring bands and their coupling $V$. That is why the dependence $\rho_{n, 0} \sim(\mathcal{G} N)^{-2 n}$ coincides with the perturbation theory limit (7.89) for the relay of weakly coupled isolated levels. When the interaction $V$ exceeds the typical distance $\delta$ between the neighboring levels of a band, the population of bands with large $n$ becomes larger. In Fig. 8.7 we show the dependence (8.68) as a function of the band number $n$ and the ratio $k / N$ for two different cases (a) $\mathcal{G} N \ll 1$ and (b) $\mathcal{G} N \gg 1$. In order to illustrate the population distribution over the bands we scale the band population by an exponential factor. The most interesting feature of the distribution is the spikes of the population near the position $k \delta \sim \mathcal{V}$ corresponding to the cooperative "8.68Rabi frequency" $\mathcal{V}$.

### 8.2 Random Walks and Coherent Behavior

While considering the behavior of the level-band systems in Chap. 3, we have seen that for time intervals exceeding the inverse energy distance between neighboring levels, the recurrences and revivals bring the population partially back to the initial state. This phenomenon results in an incomplete decay of the levels. For infinite systems such as relays, the density of states approaches infinity, thus making the estimate of the revival time infinite. However, from the results discussed in Chap. 7, we know that even infinite relays with irregularly positioned energy levels manifest the localization phenomenon, which implies that an initially populated level does not experience complete decay. This means that in the general case, the estimate of the time when the revivals and recurrences begin to play the dominating role is governed by other parameters, and therefore must be improved. We have already encountered an example of incomplete decay to an infinitely dense spectrum when considering, in Sect. 3.4, a band with extremely irregular behavior of the matrix elements, coupling states of the band to the isolated level. The incomplete decay may occur when the density of states with vanishing coupling dominates in the total state density. In this section we will trace more attentively the relation between the random walks in an infinite complex multilevel system and the quantum recurrences resulting either in non-exponential or even incomplete decay of the initially populated state. Prior to this, we consider manifestation of the returns of a random walk on the spectral properties of a quantum system.

### 8.2.1 Level Decay to a Band of Random Walks

Let us consider two coupled levels, one of which interacts with a band. Such a system has already been considered in the context of the Fano problem shown in Fig. 3.8(a). However, now we assume that only one of these two levels is coupled to a band, as was the case in Fig. 2.24(b) and that the
transition amplitudes coupling this level with the energy eigenstates of the band do not depend smoothly on the state energy, resembling the case shown in Fig. 2.5(c), and corresponding to the case considered in p . 108 and to (3.110). The Shrödinger equation for the level scheme has the same structure

$$
\begin{align*}
& i \dot{\psi}_{0}=\Delta_{0} \psi_{0}+V_{01} \psi_{1} \\
& i \dot{\psi}_{1}=\Delta_{1} \psi_{1}+V_{10} \psi_{0}+\sum_{n} V_{1 n} \psi_{n} \\
& i \dot{\psi}_{n}=\Delta_{n} \psi_{n}+V_{n 1} \psi_{1} \tag{8.69}
\end{align*}
$$

as (3.48) for the Fano problem, although the coupling strength $V_{1 n}$ now changes strongly, by orders of magnitude, for the neighboring levels, and therefore cannot be factored out from the sum. Here $\psi_{0}, \psi_{1}$ and $\psi_{n}$ are amplitudes of the states $|0\rangle,|1\rangle$ and $|n\rangle$ respectively, and $\Delta_{n}$ is the detuning of the state $|n\rangle$ in the band from the state $|1\rangle$.

In this section we do not specify the matrix elements of the coupling $V_{1 n}$, but make use of the phenomenological diffusion model. To this end we perform Fourier transformation of (8.69), express $\psi_{n}$ in terms of $\psi_{1}$ with the help of the third equation, and substitute the result into the second equation. This yields

$$
\begin{align*}
& \varepsilon \psi_{0}(\varepsilon)=V_{01} \psi_{1}(\varepsilon)+\Delta_{0} \psi_{0}(\varepsilon)+i \\
& \varepsilon \psi_{1}(\varepsilon)=\Delta_{1} \psi_{1}+V_{10} \psi_{0}(\varepsilon)+\sum_{n} \frac{V_{1 n} V_{n 1}}{\varepsilon-\Delta_{n}+i 0} \psi_{1}(\varepsilon) \tag{8.70}
\end{align*}
$$

and hence

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{i}{\varepsilon-\Delta_{0}-V_{01} V_{10}\left(\varepsilon-\Delta_{1}-\sum_{n} \frac{V_{1 n} V_{n 1}}{\varepsilon-\Delta_{n}+i 0}\right)^{-1}} \tag{8.71}
\end{equation*}
$$

The sum

$$
\begin{equation*}
G_{b}(\varepsilon)=\sum_{n} \frac{V_{1 n} V_{n 1}}{\varepsilon-\Delta_{n}+i 0} \tag{8.72}
\end{equation*}
$$

in (8.71) has a clear physical meaning: up to the factor $\left\langle V^{2}\right\rangle$ it coincides with the Fourier transform of the probability amplitude to return back to the state $|b\rangle=\sum V_{1 n} / V|n\rangle$ at time $t$, once the system leaves this state for the band at time $t=0$. Here $V=\left\langle V^{2}\right\rangle^{1 / 2}$.

In this section we take the Fourier transform $G_{b}(t)$ of $G_{b}(\varepsilon)$, proportional to the square root of the probability

$$
\begin{equation*}
\rho_{b}(t)=\frac{\exp \left\{-a^{2} / 4 D t\right\}}{(2 \pi D t)^{-d / 2}} \tag{8.73}
\end{equation*}
$$

to find the system in the $|b\rangle$ state, assuming that this probability corresponds to the classical diffusion in a $d$-dimensional space. This yields

$$
\begin{equation*}
G_{b}(\varepsilon)=-i\left\langle V^{2}\right\rangle e^{i \theta} \int_{0}^{\infty}(2 \pi D t)^{-d / 4} \exp \left\{i \varepsilon t-a^{2} / 8 D t\right\} d t \tag{8.74}
\end{equation*}
$$

where $e^{i \theta}$ is an arbitrary phase factor, and hence

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{i}{\varepsilon-\Delta_{0}-\left|V_{01}\right|^{2}\left[\varepsilon-\Delta_{1}+i V^{2} e^{i \theta} \int_{0}^{\infty}(2 \pi D t)^{-d / 4} e^{i \varepsilon t-a^{2} / 8 D t} d t\right]^{-1}} \tag{8.75}
\end{equation*}
$$

Here a finite parameter $a$ ensures the convergence of the integral (8.74) and may be interpreted as the typical size of a domain occupied by the state $|b\rangle$ in the space where the diffusion occurs. We assume that only this domain is coupled directly to the level $|1\rangle$ and take its size to be small. Note that for $a \rightarrow 0$ the integrand in (8.75) behaves as $\varepsilon^{d / 4-1}$ for $d \neq 4$ and as $\ln (\varepsilon)$ for $d=4$, and (8.75) takes the form

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{i}{\varepsilon-\Delta_{0}-\left|V_{01}\right|^{2}\left[\varepsilon-\Delta_{1}-e^{i \theta} V^{2} \varepsilon^{d / 4-1}(2 \pi i D)^{-d / 4} \Gamma\left(1-\frac{d}{4}\right)\right]^{-1}} \tag{8.76}
\end{equation*}
$$

Equation (8.76) has two simple poles in the complex plane corresponding to the exponentially decreasing contributions to the population $\rho_{0}(t)=$ $\left|\psi_{0}(t)\right|^{2}$ and a branching point at $\varepsilon=0$ which yields a power dependence dominating in the limit of long times. We concentrate on the last contribution which is the most important one and consider two different cases $0<d<4$ and $4<d<8$ corresponding to the different asymptotic behavior of the Green's functions at long times. For $d>4$ one finds that $G_{b}(\varepsilon) \rightarrow 0$ at $\varepsilon \rightarrow 0$, and the Taylor expansion over $G_{b}(\varepsilon)$ yields the first non-vanishing contribution

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{i\left|V_{01}\right|^{2} e^{i \theta} V^{2}}{\left(\Delta_{0} \Delta_{1}-\left|V_{01}\right|^{2}\right)^{2}} \varepsilon^{d / 4-1}(2 \pi i D)^{-d / 4} \Gamma\left(1-\frac{d}{4}\right) \tag{8.77}
\end{equation*}
$$

provided $\Delta_{0} \Delta_{1} \neq V_{01} V_{10}$. For $\Delta_{0} \Delta_{1}=V_{01} V_{10}$ the expansion reads

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{-i \Delta_{1} e^{-i \theta}}{\Delta_{0} \varepsilon^{d / 4-1} V^{2}(2 \pi i D)^{-d / 4} \Gamma\left(1-\frac{d}{4}\right)} \tag{8.78}
\end{equation*}
$$

For $d<4$ the Green's function $G_{b}(\varepsilon) \rightarrow \infty$ for $\varepsilon \rightarrow 0$, and for $\Delta_{0} \neq 0$ the Taylor expansion over $1 / G_{b}(\varepsilon)$ results in

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{\varepsilon^{1-d / 4}}{i \Delta_{0}^{2}} \frac{e^{-i \theta}\left|V_{01}\right|^{2}}{V^{2}(2 \pi i D)^{-d / 4} \Gamma\left(1-\frac{d}{4}\right)} \tag{8.79}
\end{equation*}
$$

whereas for $\Delta_{0}=0$ it yields

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{e^{i \theta} \varepsilon^{d / 4-1} V^{2}(2 \pi i D)^{-d / 4} \Gamma\left(1-\frac{d}{4}\right)}{\left|V_{01}\right|^{2}} \tag{8.80}
\end{equation*}
$$

For $d>4$ after inverse Fourier transformation (8.77) and (8.78) yield

$$
\begin{equation*}
\rho_{0}(t) \simeq \frac{V^{4}\left|V_{01}\right|^{4} \pi^{2}}{(2 \pi D t)^{d / 2}\left(\Delta_{0} \Delta_{1}-\left|V_{01}\right|^{2}\right)^{4}} \tag{8.81}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{0}(t) \simeq \frac{(2 \pi D t)^{d / 2} \Delta_{1}^{2}(1-d / 4)^{2} \sin ^{2}(\pi d / 4)}{V^{4} t^{4} \Delta_{0}^{2}} \tag{8.82}
\end{equation*}
$$

respectively. This means that the population of the state $|0\rangle$ manifests the time behavior suggested by the diffusion. However, at resonance, when the detuning of the ground state shifted by the interaction with the state|1〉 vanishes, that is $\Delta_{0}-\left|V_{01}\right|^{2} / \Delta_{1}=0$, the decay gets slower by the factor $\sim t^{d-4}$.

For $d<4$, equations (8.79)-(8.80) result in

$$
\begin{equation*}
\rho_{0}(t) \simeq \frac{(2 \pi D t)^{d / 2}\left|V_{01}\right|^{4}(1-d / 4)^{2} \sin ^{2}(\pi d / 4)}{V^{4} t^{4} \Delta_{0}^{4}} \tag{8.83}
\end{equation*}
$$

and

$$
\begin{equation*}
\rho_{0}(t) \simeq \frac{\pi^{2} V^{4}}{(2 \pi D t)^{d / 2}\left|V_{01}\right|^{4}} \tag{8.84}
\end{equation*}
$$

In this regime, the behavior is similar - the slowest decay occurs in the resonance $\Delta_{0}=0$, but corresponds now to diffusive behavior. The off-resonant decay is faster by the factor $\sim t^{d-4}$. We note that the presence of the state $|1\rangle$ no longer result in a shift of the resonance between the band and the state $|0\rangle$, since the state $|1\rangle$ is strongly broadened due to the interaction with the continuum.

The crucial role of returns of the population amplitudes (or the population amplitude wavepackets) in the course of a random walk becomes much more clear in another particular case related to the level-band problem. Let us assume that a probability amplitude is injected into the state $|0\rangle$ with the rate $\Pi(\varepsilon)$, similar to the case shown in Fig. 3.3 for a level-band system. Equation (8.76) then adopts the form

$$
\begin{equation*}
\psi_{0}(\varepsilon)=\frac{i \Pi(\varepsilon)}{\varepsilon-\Delta_{0}-\left|V_{01}\right|^{2}\left[\varepsilon-\Delta_{1}-e^{i \theta} V^{2} \varepsilon^{d / 4-1}(2 \pi i D)^{-d / 4} \Gamma\left(1-\frac{d}{4}\right)\right]^{-1}} \tag{8.85}
\end{equation*}
$$

From (8.85) one finds the Fourier transform $\rho(t)$

$$
\begin{align*}
\rho_{0}(t) & =\frac{1}{2 \pi} \int_{-\infty}^{\infty} \psi_{0}(\varepsilon) \psi_{0}^{*}(\xi) e^{-i(\varepsilon-\xi) t} d \varepsilon d \xi  \tag{8.86}\\
& =\frac{1}{2 \pi} \int \frac{\Pi(\varepsilon) \Pi(\xi)}{\left[\varepsilon-\Delta_{0}-\frac{\left|V_{01}\right|^{2}(\varepsilon)}{\varepsilon-\Delta_{1}-G_{b}(\varepsilon)}\right]\left[\xi-\Delta_{0}-\frac{\left|V_{01}\right|^{2}}{\xi-\Delta_{1}-G_{b}^{*}(\xi)}\right]} e^{-i(\varepsilon-\xi) t} d \varepsilon d \xi
\end{align*}
$$

of the probability to be at the state $|0\rangle$ at time $t$. Assuming the injected population amplitude $\Pi(t)$ to be a random function, with a $\delta$-like time autocorrelation that results in $\langle\Pi(\varepsilon) \Pi(\xi)\rangle \sim \delta(\varepsilon-\xi)$, one finds a steady state population


Fig. 8.8. Population flux $\Phi$ to a band corresponding to a multidimensional random walk. The flux $\Phi$ for $V_{01}=1.3$ as a function of the detunings $\Delta_{0}$ and $\Delta_{1}$, for different dimensions $d=3$ (a) and $d=5$ (b). The flux $\Phi$ for $\Delta_{1}=1.3$ as a function of $\Delta_{0}$ and $V_{01}$ for $d=3$ (c) and $d=5$ (d).

$$
\begin{equation*}
\rho_{0} \sim \int \frac{d \varepsilon}{\left[\varepsilon-\Delta_{0}-\frac{\left|V_{01}\right|^{2}}{\varepsilon-\Delta_{1}-G_{b}(\varepsilon)}\right]\left[\varepsilon-\Delta_{0}-\frac{\left|V_{01}\right|^{2}}{\varepsilon-\Delta_{1}-G_{b}^{*}(\varepsilon)}\right]} . \tag{8.87}
\end{equation*}
$$

If we assume that the amplitude injection is saturated, such that population of the state $|0\rangle$ remains of the order of its maximum possible value $\rho_{0} \simeq 1$ then the flux $\Phi$ of the population to the continuum corresponding to the random walk is given by the inverse value of the integral (8.87).

$$
\begin{equation*}
\Phi \sim 1 / \int \frac{d \varepsilon}{\left[\varepsilon-\Delta_{0}-\frac{\left|V_{01}\right|^{2}}{\varepsilon-\Delta_{1}-G_{b}(\varepsilon)}\right]\left[\varepsilon-\Delta_{0}-\frac{\left|V_{01}\right|^{2}}{\varepsilon-\Delta_{1}-G_{b}^{*}(\varepsilon)}\right]} . \tag{8.88}
\end{equation*}
$$

In Fig. 8.8 we show this flux $\Phi$ as a function of the detunings $\Delta_{0}$ and $\Delta_{1}$, and as a function of $\Delta_{0}$ and $V_{01}$ for different dimensions $d$ of the random walk. Diffusive returns of the population amplitude play an important role inhibiting the decay at small detuning $\Delta_{0}$ and result in a depth in the center of the line at $d<4$. In other words, interference of the population decay with these recurrences strongly affects the line shapes of the resonances. For $d>4$ returns of the population amplitude do not occur, and the line shape resembles a Lorentzian profile.

### 8.2.2 Interference of Random Returns at Long Times. General Consideration

We now consider the behavior of quantum systems for asymptotically long times. For a complex system, this behavior might be strongly affected by the important role played by self-intersections in the diagrams. Under certain conditions, the main contribution to the perturbation series give the diagrams that return many times at the same quantum levels, and allow for the selfinterference of the wavepackets. The most important characteristic which determines the behavior of a complex quantum system is therefore closely related to a property of the corresponding classical random walk, given by the probability for a randomly moving particle to return back to the points which it has previously visited. In the quantum limit, this corresponds to multiple self-intersections of the trajectories, which results in strong interference of many essentially different quantum paths.

The direct approach can be based on the combinatorial calculation of the numbers of different self-intersecting trajectories and the numbers of different ways to pass along these trajectories. The latter is well known in mathematics as the number of Euler trails on the oriented graphs. An explicit formula exists for such a number, as a function of the self-intersecting diagram topology, and the number of paths entering each self-intersection. However the calculations based on this intuitively clear idea are usually rather cumbersome, and therefore we demonstrate this approach in the next section, for a simpler example. Here we suggest a less instructive consideration based on the integrals over Grassmann variables, that have already been introduced in the previous chapter. Though it hides to a certain extent the intuitively meaningful images associated with the diagrams of interfering trajectories, it is most practical technically, giving more compact intermediate expressions.

## Level-Band-Random Walk Continuum Model

We consider the system shown in Fig. 8.9(a) consisting of a non-degenerate


Fig. 8.9. The type of random walk affects the level decay at long times. (a) A nondegenerate level $|0\rangle$ is coupled directly by interactions $V_{0 n}$ of a random strength to a discrete set of quantum states randomly distributed in energy that are coupled among themselves via the states $|\kappa\rangle$ of an irregular continuum. The Fourier transforms of the decay rate $G(\varepsilon)$ to the continuum and the correlation function $\lambda^{2}(\varepsilon, \xi)$ of the return amplitudes are the main characteristics of the process. (b) The unlikely $(\sim \exp (-\Delta g))$, long decaying realizations of the band, with a large gap $\Delta$ in the resonance position, determine the long-time behavior for a non-returnable random walk. (c) Two types of returnable random walks. The exponential wavepackets leaving the band states acquire, after return from the continuum, either zero or a finite autocorrelation time.
level $|0\rangle$ directly coupled to a band of levels $|n\rangle$. The latter are coupled among themselves by an indirect interaction which occurs via a very dense or even a continuous set of states $|\kappa\rangle$ that represents a random walk with certain statistical properties. The Schrödinger equation for the system initially in the state $|0\rangle$ reads

$$
\begin{align*}
i \dot{\psi}_{0} & =V_{0 n} \psi_{n}+i \delta(t) \\
i \dot{\psi}_{n} & =\Delta_{n} \psi_{n}+V_{n 0} \psi_{0}+\sum_{\kappa} V_{n \kappa} \psi_{\kappa} \\
i \dot{\psi}_{\kappa} & =\Delta_{\kappa} \psi_{\kappa}+\sum_{n} V_{\kappa n} \psi_{n} \tag{8.89}
\end{align*}
$$

and adopts the form

$$
\begin{align*}
& \varepsilon \psi_{0}=V_{0 n} \psi_{n}+i \\
& \varepsilon \psi_{n}=\Delta_{n} \psi_{n}+V_{n 0} \psi_{0}+\sum_{\kappa} V_{n \kappa} \psi_{\kappa} \\
& \varepsilon \psi_{\kappa}=\Delta_{\kappa} \psi_{\kappa}+\sum_{n} V_{\kappa n} \psi_{n} \tag{8.90}
\end{align*}
$$

after the Fourier transformation. We assume a uniform distribution $\Delta_{n}$ of the level positions in the entire interval from $-\infty$ to $\infty$ and a Gaussian distribution of $e^{-\left|V^{2}\right|}$ complex valued couplings $V_{0 n}$ with unit variance. All energy variables are therefore scaled by the mean square coupling $V=\left\langle V^{2}\right\rangle^{1 / 2}$ whereas the mean density of states $g$ is understood as a large dimensionless parameter $g V$, representing the expectation number of resonances, that have been considered in Sect.3.4.4.

Introducing in complete analogy with (8.72) the function

$$
\begin{equation*}
G_{n n^{\prime}}(\varepsilon)=\sum_{\kappa} \frac{V_{n \kappa} V_{\kappa n^{\prime}}}{\varepsilon-\Delta_{\kappa}} \tag{8.91}
\end{equation*}
$$

we exclude from the consideration the levels of the band representing the random walk and the Schrödinger equation adopts the form

$$
\begin{align*}
& \varepsilon \psi_{0}=V_{0 n} \psi_{n}+i \\
& \varepsilon \psi_{n}=\Delta_{n} \psi_{n}+V_{n 0} \psi_{0}+\sum_{n^{\prime}} G_{n n^{\prime}}(\varepsilon) \psi_{n^{\prime}} \tag{8.92}
\end{align*}
$$

A similar equation holds for the Fourier transforms of the complex conjugate amplitudes $\psi_{n}(\xi)$ and the corresponding functions $G_{n n^{\prime}}(\xi)$.

The moments $\left\langle G_{n n^{\prime}}^{k}(\varepsilon)\right\rangle$ and the correlations $\left\langle G_{n n^{\prime}}^{k}(\varepsilon) G_{m m^{\prime}}^{k^{\prime}}(\xi)\right\rangle$ of the functions (8.71) contain information about all of the statistical properties of the random walk. Later on we will discuss in detail the particular properties governing the asymptotic behavior of the quantum systems in the long-time limit. Here we just specify several general hypotheses employed. We assume that the random walk is completely defined by the first $G(\varepsilon)=\left\langle G_{n n^{\prime}}(\varepsilon)\right\rangle$ and the second $\lambda^{2}(\varepsilon, \xi)=\left\langle G_{n n^{\prime}}(\varepsilon) G_{m m^{\prime}}(\xi)\right\rangle$ moments. In other words, the Green's functions are well-defined, have zero dispersion in absolute values, and vary randomly in phase. Therefore the non-vanishing higher moments are explicitly given by the two first moments via the relations.

$$
\begin{align*}
\left\langle G_{n n^{\prime}}^{k}(\varepsilon) G_{m m^{\prime}}^{k^{\prime}}(\xi)\right\rangle & =\delta_{k k^{\prime}} \delta_{n m} \delta_{n^{\prime} m^{\prime}}\left[\lambda^{2}(\varepsilon, \xi)\right]^{k} \\
\left\langle G_{n n^{\prime}}^{k}(\varepsilon)\right\rangle & =\delta_{n n^{\prime}}[G(\varepsilon)]^{k} \tag{8.93}
\end{align*}
$$

## Effective Hamiltonians and Averaging over the Random Walk

Once the amplitudes of the states $|\kappa\rangle$ are excluded, the dynamics of the remaining system consisting of the level $|0\rangle$ and the band of the states $|n\rangle$ satisfy the Schrödinger equation (8.92) with the effective Hamiltonian

$$
\widehat{H}(\varepsilon)=\left(\begin{array}{ccccc}
0 & V_{01} & V_{02} & \ldots & V_{0 N}  \tag{8.94}\\
V_{10} & \Delta_{1}+G_{11}(\varepsilon) & G_{12}(\varepsilon) & \ldots & G_{1 N}(\varepsilon) \\
V_{20} & G_{21}(\varepsilon) & \Delta_{2}+G_{22}(\varepsilon) & \ldots & G_{2 N}(\varepsilon) \\
\ldots & \ldots & \ldots & \ldots & \ldots \\
V_{0 N} & G_{N 1}(\varepsilon) & G_{N 2}(\varepsilon) & \ldots & \Delta_{N}+G_{N N}(\varepsilon)
\end{array}\right) .
$$

The Fourier transform of the population $\rho_{0}(\varepsilon, \xi)$ of the state $|0\rangle$ reads

$$
\rho_{0}=\frac{\left|\begin{array}{cccc}
\Delta_{1 \varepsilon} & G_{12}(\varepsilon) & \ldots & G_{1 N}(\varepsilon)  \tag{8.95}\\
G_{21}(\varepsilon) & \Delta_{2 \varepsilon} & \ldots & G_{2 N}(\varepsilon) \\
\ldots & \ldots & \ldots & \ldots \\
G_{N 1}(\varepsilon) & G_{N 2}(\varepsilon) & \ldots & \Delta_{N \varepsilon}
\end{array}\right|\left|\begin{array}{ccccc}
\Delta_{1 \xi} & G_{12}(\xi) & \ldots & G_{1 N}(\xi) \\
G_{21}(\xi) & \Delta_{2 \xi} & \ldots & G_{2 N}(\xi) \\
\ldots & \ldots & \ldots & \ldots \\
G_{N 1}(\xi) & G_{N 2}(\xi) & \ldots & \Delta_{N \xi}
\end{array}\right|}{\left|\begin{array}{ccccc}
-\varepsilon & V_{01} & \ldots & V_{0 N} \\
V_{10} & \Delta_{1 \varepsilon} & \ldots & G_{1 N}(\varepsilon) \\
\ldots & \ldots & \ldots & \ldots \\
V_{0 N} & G_{N 1}(\varepsilon) & \ldots & \Delta_{N \varepsilon}
\end{array}\right|\left|\begin{array}{cccc}
-\xi & V_{01} & \ldots & V_{0 N} \\
V_{10} & \Delta_{1 \xi} & \ldots & G_{1 N}(\xi) \\
\ldots & \ldots & \ldots & \ldots \\
V_{0 N} & G_{N 1}(\xi) & \ldots & \Delta_{N \xi}
\end{array}\right|}
$$

as it directly follows from the solution of (8.92). Here the replacements $\Delta_{n}+$ $G_{n n}(\xi)-\xi \rightarrow \Delta_{n \xi}$ and $\Delta_{n}+G_{n n}(\varepsilon)-\varepsilon \rightarrow \Delta_{n \varepsilon}$ are performed to shorten the notation.

We have to find the average of the population over all possible realizations of the couplings $V_{0 n}$ and the positions $\Delta_{n}$ of the levels which we assume to be statistically independent, that is, following Poissinian statistics. In the representation (8.95) this is impossible, since $\rho_{0}(\varepsilon, \xi)$ depends on complex combinations of $V_{0 n}$ and $\Delta_{n}$. A possible approach should rely on factorization of the expression for the population in such a way that each cofactor corresponds to a different $|n\rangle$ and is distributed in the same way. To this end, the ratios of the determinants can be represented as Gaussian integrals over complex numbers $x$ and the Grassmann variables $\chi$ with the help of (7.133), (7.136), and (7.180). This yields

$$
\begin{align*}
\rho_{0}(\varepsilon, \xi)= & \int e^{-\sum_{n, n^{\prime}=1}^{N}\left[\widehat{H}_{n n^{\prime}}^{\prime}(\varepsilon) \bar{\chi}_{n} \chi_{n^{\prime}}+\widehat{H}_{n n^{\prime}}^{\prime}(\xi) \bar{\chi}_{n}^{\prime} \chi_{n^{\prime}}^{\prime}\right]} \prod_{n=1}^{N} d \bar{\chi}_{n} d \chi_{n^{\prime}} d \bar{\chi}_{n}^{\prime} d \chi_{n^{\prime}}^{\prime} \\
& \times \int e^{-\sum_{n, n^{\prime}=0}^{N}\left[\widehat{H}_{n n^{\prime}}^{\prime}(\varepsilon) \bar{x}_{n} x_{n^{\prime}}+\widehat{H}_{n n^{\prime}}^{\prime}(\xi) \bar{x}_{n}^{\prime} x_{n^{\prime}}^{\prime}\right]} \prod_{n=0}^{N} d \bar{x}_{n} d x_{n} d \bar{x}_{n}^{\prime} d x_{n}^{\prime}, \tag{8.96}
\end{align*}
$$

where we denote $\widehat{H}_{n n^{\prime}}(\varepsilon)-\varepsilon=\widehat{H}_{n n^{\prime}}^{\prime}(\varepsilon)$ and $\widehat{H}_{n n^{\prime}}(\xi)-\xi=\widehat{H}_{n n^{\prime}}^{\prime}(\xi)$ to shorten the notation. For the same reason, the normalization $1 / \pi$ of Gaussian integrals over $d \bar{x}_{n} d x_{n}$ is implicit.

We first make use of the statistical properties of the random walk. We cast the right-hand side of (8.96)

$$
\begin{align*}
& \int \prod_{n=0}^{N} d \bar{x}_{n} d x_{n^{\prime}} d \bar{x}_{n}^{\prime} d x_{n^{\prime}}^{\prime} \prod_{n=1}^{N} d \bar{\chi}_{n} d \chi_{n^{\prime}} d \bar{\chi}_{n}^{\prime} d \chi_{n^{\prime}}^{\prime} \exp \left\{i \varepsilon \bar{x}_{0} x_{0}-i \xi \bar{x}_{0}^{\prime} x_{0}^{\prime}\right. \\
& -i \sum_{n=1}^{N} \Delta_{n \varepsilon}\left(\bar{\chi}_{n} \chi_{n}+\bar{x}_{n} x_{n}\right)+i \sum_{n=1}^{N} \Delta_{n \xi}\left(\bar{\chi}_{n}^{\prime} \chi_{n}^{\prime}+\bar{x}_{n}^{\prime} x_{n}^{\prime}\right) \\
& +i \sum_{n=1}^{N} V_{0 n}\left(\bar{x}_{0} x_{n}-\bar{x}_{n}^{\prime} x_{0}^{\prime}\right)-i \sum_{n=1}^{N} V_{0 n}^{*}\left(\bar{x}_{0}^{\prime} x_{n}^{\prime}-\bar{x}_{n} x_{0}\right)  \tag{8.97}\\
& \left.+\sum_{n, n^{\prime}=1}^{N}\left[i G_{n^{\prime} n}(\varepsilon)\left(\bar{\chi}_{n^{\prime}} \chi_{n}+\bar{x}_{n^{\prime}} x_{n}\right)-i G_{n^{\prime} n}(\xi)\left(\bar{\chi}_{n^{\prime}}^{\prime} \chi_{n}^{\prime}+\bar{x}_{n^{\prime}}^{\prime} x_{n}^{\prime}\right)\right]\right\}
\end{align*}
$$

in a Taylor series over $G$, make use of (8.93) performing the average over the random walks, and employ the integral representation for the factor $1 / k$ ! and write

$$
\begin{equation*}
\frac{\left\langle G_{n n^{\prime}}^{k}(\varepsilon) G_{m m^{\prime}}^{k}(\xi)\right\rangle}{k!k!}=\delta_{n m} \delta_{n^{\prime} m^{\prime}} \frac{1}{k!} \int_{C} \frac{e^{-\tau}}{2 \pi i \tau}\left[\frac{\lambda^{2}(\varepsilon, \xi)}{-\tau}\right]^{k} d \tau \tag{8.98}
\end{equation*}
$$

where the contour $C$ circumvents the point $\tau=0$ in the positive direction starting and ending at $\tau=\infty$. Summing the Taylor series back we arrive at

$$
\begin{align*}
\left\langle\rho_{0}(\varepsilon, \xi)\right\rangle= & \int \frac{e^{-\tau}}{2 \pi i \tau} d \tau \prod_{n=0}^{N} d \bar{x}_{n} d x_{n^{\prime}} d \bar{x}_{n}^{\prime} d x_{n^{\prime}}^{\prime} \prod_{n=1}^{N} d \bar{\chi}_{n} d \chi_{n^{\prime}} d \bar{\chi}_{n}^{\prime} d \chi_{n^{\prime}}^{\prime} \\
& \times \exp \left\{-i \sum_{n=1}^{N} \Delta_{n \varepsilon}\left(\bar{\chi}_{n} \chi_{n}+\bar{x}_{n} x_{n}\right)+i \sum_{n=1}^{N} \Delta_{n \xi}\left(\bar{\chi}_{n}^{\prime} \chi_{n}^{\prime}+\bar{x}_{n}^{\prime} x_{n}^{\prime}\right)\right. \\
& +i \sum_{n=1}^{N} V_{0 n}\left(\bar{x}_{0} x_{n}-\bar{x}_{n}^{\prime} x_{0}^{\prime}\right)-i \sum_{n=1}^{N} V_{0 n}^{*}\left(\bar{x}_{0}^{\prime} x_{n}^{\prime}-\bar{x}_{n} x_{0}\right)+i \varepsilon \bar{x}_{0} x_{0} \\
& \left.-i \xi \bar{x}_{n}^{\prime} x_{n}^{\prime}-\frac{\lambda^{2}}{\tau} \sum_{n, n^{\prime}=1}^{N}\left(\bar{\chi}_{n^{\prime}} \chi_{n}+\bar{x}_{n^{\prime}} x_{n}\right)\left(\bar{\chi}_{n^{\prime}}^{\prime} \chi_{n}^{\prime}+\bar{x}_{n^{\prime}}^{\prime} x_{n}^{\prime}\right)\right\} \tag{8.99}
\end{align*}
$$

which now has to be averaged only over the ensemble of the band levels.

## Factorization of the Population

Now our aim is to represent (8.99) in factorized form. To this end, in the exponent of (8.99), we multiply factors proportional to $\lambda^{2}$ and regroup them differently

$$
\begin{align*}
& \exp \left[-\frac{\lambda^{2}}{\tau} \sum_{n, n^{\prime}=1}^{N}\left(\bar{\chi}_{n^{\prime}} \chi_{n}+\bar{x}_{n^{\prime}} x_{n}\right)\left(\bar{\chi}_{n^{\prime}}^{\prime} \chi_{n}^{\prime}+\bar{x}_{n^{\prime}}^{\prime} x_{n}^{\prime}\right)\right] \\
= & \exp \left[-\frac{\lambda^{2}}{\tau}\left(S_{\overline{\chi \chi^{\prime}}} S_{\chi^{\prime} \chi}+\sigma_{\overline{\chi x}} \sigma_{x \chi^{\prime}}+\sigma_{x^{\prime} \bar{\chi}} \sigma_{\chi \bar{x}^{\prime}}+S_{\bar{x} x^{\prime}} S_{x \bar{x}^{\prime}}\right)\right] \tag{8.100}
\end{align*}
$$

with the help of the combinations

$$
\begin{align*}
& \sum_{n=1}^{N} \bar{x}_{n} x_{n}^{\prime}=S_{\bar{x} x^{\prime}} ; \sum_{n=1}^{N} x_{n} \bar{x}_{n}^{\prime}=S_{x \bar{x}^{\prime}} ; \sum_{n=1}^{N} \bar{\chi}_{n} \bar{\chi}_{n}^{\prime}=S_{\overline{\chi \chi^{\prime}}} ; \sum_{n=1}^{N} \chi_{n}^{\prime} \chi_{n}=S_{\chi^{\prime} \chi} \\
& \sum_{n=1}^{N} \bar{\chi}_{n}^{\prime} \bar{x}_{n}=\sigma_{\overline{\chi x}} ; \sum_{n=1}^{N} x_{n} \chi_{n}^{\prime}=\sigma_{x \chi^{\prime}} ; \sum_{n=1}^{N} x_{n}^{\prime} \bar{\chi}_{n}=\sigma_{x^{\prime} \bar{\chi}} ; \sum_{n=1}^{N} \chi_{n} \bar{x}_{n}^{\prime}=\sigma_{\chi \bar{x}^{\prime}} \tag{8.101}
\end{align*}
$$

that do not contain double sums. The combinations $S$ contain complex numbers and bi-Grassmann variables, that from the viewpoint of the commutation relations are equivalent to complex numbers, while the combinations $\sigma$ are equivalent to Grassmann variables. We note the following relations for Gaussian integrals over complex conjugated variables $z=x+i y$ and $\bar{z}=x-i y$

$$
\begin{align*}
\int_{-\infty}^{\infty} \frac{d z d \bar{z}}{\pi} \exp (i a z+i b \bar{z}-z \bar{z}) & =\int_{-\infty}^{\infty} \frac{d x d y}{\pi} e^{i(a+b) x+(b-a) y-x^{2}-y^{2}} \\
& =e^{(b-a)^{2} / 4-(a+b)^{2} / 4}=e^{-a b} \tag{8.102}
\end{align*}
$$

are valid for any numbers $a$ and $b$. Therefore the factorization of the binary combinations of the sums $S$ in the exponent can be performed with the help of the integrals over two pairs of complex conjugated numbers - the pair $f$ and $\bar{f}$ for $S_{\bar{x} x^{\prime}} S_{x \bar{x}^{\prime}}$ and the pair $w$ and $\bar{w}$ for $S_{\overline{\chi \chi^{\prime}}} S_{\chi^{\prime} \chi}$ :

$$
\begin{align*}
& \exp \left[-\frac{\lambda^{2}}{\tau}\left(S_{\overline{\chi \chi^{\prime}}} S_{\chi^{\prime} \chi}+S_{\bar{x} x^{\prime}} S_{x \bar{x}^{\prime}}\right)\right]=\int d w d \bar{w} d f d \bar{f} \\
& \exp \left[\frac{i \lambda}{\sqrt{\tau}}\left(S_{\overline{\chi \chi}} w+\bar{w} S_{\chi^{\prime} \chi}+S_{\bar{x} x^{\prime}} f+\bar{f} S_{x \bar{x}^{\prime}}\right)-w \bar{w}-f \bar{f}\right] \tag{8.103}
\end{align*}
$$

where the normalization factors $1 / \pi$ for $d w d \bar{w}$ and for $d f d \bar{f}$ are implicit. For Grassmann combinations $\sigma$ one finds an analog of the relation (8.102)

$$
\begin{equation*}
\int \exp [-\bar{\alpha} \alpha-\bar{\sigma} \alpha-\bar{\alpha} \sigma] d \bar{\alpha} d \alpha=\exp \bar{\sigma} \sigma \tag{8.104}
\end{equation*}
$$

With the help of (8.102) and (8.104) one obtains only a linear combination of sums (8.101) in the exponent, and therefore the population adopts the desired factorized form

$$
\begin{align*}
& \left\langle\rho_{0}(\varepsilon, \xi)\right\rangle=\int \frac{d \tau}{2 \pi i \tau} d w d \bar{w} d f d \bar{f} d \bar{\alpha} d \alpha d \bar{\beta} d \beta e^{-\tau-w \bar{w}-f \bar{f}-\bar{\alpha} \alpha-\bar{\beta} \beta} \\
& \left\{\prod_{n=0}^{N} \int d \bar{x}_{n} d x_{n} d \bar{x}_{n}^{\prime} d x_{n}^{\prime} d \bar{\chi}_{n} d \chi_{n} d \bar{\chi}_{n}^{\prime} d \chi_{n}^{\prime}\right\} e^{i \varepsilon \bar{x}_{0} x_{0}-i \xi \bar{x}_{0}^{\prime} x_{0}^{\prime}} \prod_{n=1}^{N} e^{L_{n}} \tag{8.105}
\end{align*}
$$

where

$$
\begin{align*}
L_{n}= & i V_{0 n}\left(\bar{x}_{0} x_{n}-\bar{x}_{n}^{\prime} x_{0}^{\prime}\right)-i V_{0 n}^{*}\left(\bar{x}_{0}^{\prime} x_{n}^{\prime}-\bar{x}_{n} x_{0}\right) \\
& -i \Delta_{n \varepsilon}\left(\bar{\chi}_{n} \chi_{n}+\bar{x}_{n} x_{n}\right)+i \Delta_{n \xi}\left(\bar{\chi}_{n}^{\prime} \chi_{n}^{\prime}+\bar{x}_{n}^{\prime} x_{n}^{\prime}\right) \\
& -\frac{\lambda}{\sqrt{\tau}}\left(x_{n}^{\prime} \bar{\chi}_{n} \alpha+\bar{\alpha} \chi_{n} \bar{x}_{n}^{\prime}+\bar{\chi}_{n}^{\prime} \bar{x}_{n} \beta+\bar{\beta} x_{n} \chi_{n}^{\prime}\right) . \\
& +\frac{i \lambda}{\sqrt{\tau}}\left(\bar{\chi}_{n} \bar{\chi}_{n}^{\prime} w+\bar{w} \chi_{n}^{\prime} \chi_{n}+\bar{x}_{n} x_{n}^{\prime} f+\bar{f} x_{n} \bar{x}_{n}^{\prime}\right) . \tag{8.106}
\end{align*}
$$

The expression for the ensemble averaged population therefore takes the form of a 9 -fold integral

$$
\begin{align*}
\left\langle\rho_{0}(\varepsilon, \xi)\right\rangle= & \int \frac{e^{-\tau}}{2 \pi i \tau} d \tau d w d \bar{w} d f d \bar{f} d \bar{\alpha} d \alpha d \bar{\beta} d \beta \exp [-w \bar{w}-f \bar{f}-\bar{\alpha} \alpha-\bar{\beta} \beta] \\
& \int d \bar{x}_{0} d x_{0} d \bar{x}_{0}^{\prime} d x_{0}^{\prime} e^{i \varepsilon \bar{x}_{0} x_{0}-i \xi \bar{x}_{0}^{\prime} x_{0}^{\prime}} \prod_{n=1}^{N}\left\langle\Psi\left(\Delta_{n}, x_{0}, \ldots, \bar{\beta}\right)\right\rangle \cdot(8.10 \tag{8.107}
\end{align*}
$$

which includes the product of identical and independently distributed factors

$$
\begin{equation*}
\Psi\left(\Delta_{n}, x_{0}, \ldots, \bar{\beta}\right)=\int d \bar{x} d x d \bar{x}^{\prime} d x^{\prime} d \bar{\chi} d \chi d \bar{\chi}^{\prime} d \chi^{\prime} \exp \left(L_{n}\right) \tag{8.108}
\end{equation*}
$$

averaged over the distribution of statistically independent positions and couplings of the band levels. The index $n$ of the integration variables $x$ and $\chi$ is omitted since all the $|\kappa\rangle$ levels are statistically equivalent, and the averages of the integrals (8.108) for $n \neq 0$ are identical. The marker $n$ of the statistically independent level position $\Delta_{n}$ can be omitted as well.

## Average over the Ensemble of the Band Levels

Our next step is to perform the ensemble average of $\Psi\left(\Delta, x_{0}, \ldots, \bar{\beta}\right)$ over the positions $\Delta$ and the couplings $V$ of the levels $|n\rangle$. For the distribution $g(V, \Delta)$ of these quantities we choose the uniform Poissonian dependence on $\Delta$ in the entire interval from $-\infty$ to $\infty$ and the Gaussian distribution $e^{-\left|V^{2}\right|}$ of the complex valued $V$. The average over the ensemble with the coupling distribution $e^{-\left|V^{2}\right|}$ is straightforward - according to (8.102) the average of the $V$-dependent part immediately yields

$$
\begin{align*}
& \int e^{-\left|V^{2}\right|} \exp \left[i V\left(\bar{x}_{0} x-\bar{x}^{\prime} x_{0}^{\prime}\right)-i V^{*}\left(\bar{x}_{0}^{\prime} x^{\prime}-\bar{x} x_{0}\right)\right] d V d V^{*} \\
= & \exp \left(-\bar{x}_{0} \bar{x}_{0}^{\prime} x x^{\prime}+x_{0}^{\prime} \bar{x}_{0}^{\prime} \bar{x}^{\prime} x^{\prime}+x_{0} \bar{x}_{0} x \bar{x}-x_{0}^{\prime} x_{0} \bar{x}^{\prime} \bar{x}\right) . \tag{8.109}
\end{align*}
$$

However the average over the ensemble of uniformly distributed $\Delta$ is a more sophisticated procedure. Indeed, for large $\Delta$ the integral over the variables $x$ and $\chi$ yields a value close to unity as one can see neglecting in (8.106) the terms of $L_{n}$, independent of $\Delta$ and performing the integration

$$
\begin{equation*}
\int d \bar{x} d x d \bar{x}^{\prime} d x^{\prime} d \bar{\chi} d \chi d \bar{\chi}^{\prime} d \chi^{\prime} e^{-i \Delta_{\varepsilon}(\bar{\chi} \chi+\bar{x} x)+i \Delta_{\xi}\left(\bar{\chi}^{\prime} \chi^{\prime}+\bar{x}^{\prime} x^{\prime}\right)}=\frac{-i \Delta_{\varepsilon} i \Delta_{\xi}}{-i \Delta_{\varepsilon} i \Delta_{\xi}}=1 \tag{8.110}
\end{equation*}
$$

Therefore the ensemble average of $\Psi\left(\Delta, x_{0}, \ldots, \bar{\beta}\right)$ is a value infinitesimally close to 1 . However the product of an infinite number of such factors can be arbitrary. We have already encountered such a situation in Chap. 3, considering the decay of the level to a dense and irregular band. The proper way to overcome the problem has been demonstrated in (3.118). Employing the relation $\prod \Psi\left(\Delta_{n}\right) \simeq \exp \left\{\sum\left[\Psi\left(\Delta_{n}\right)-1\right]\right\}$ one can put the small values $\Psi-1$
into the exponent, and perform the integration of $\Psi$ over $\Delta$ there, closing the integration contour $C$ at infinity where the integral equals unity, that is $\int_{-\infty}^{\infty}[\Psi(\Delta)-1] d \Delta \rightarrow \int_{C} \Psi(\Delta) d \Delta$. For the population (8.107) this yields

$$
\begin{aligned}
\left\langle\rho_{0}(\varepsilon, \xi)\right\rangle= & \int \frac{e^{-\tau}}{\tau} d \tau d w d \bar{w} d f d \bar{f} d \bar{\alpha} d \alpha d \bar{\beta} d \beta e^{-w \bar{w}-f \bar{f}-\bar{\alpha} \alpha-\bar{\beta} \beta} \\
& \int d \bar{x}_{0} d x_{0} d \bar{x}_{0}^{\prime} d x_{0}^{\prime} e^{i \varepsilon \bar{x}_{0} x_{0}-i \xi \bar{x}_{0}^{\prime} x_{0}^{\prime}} \exp \left(\int_{C} \mathbf{g} \Psi d V d V^{*} d \Delta\right)-
\end{aligned}
$$

where $\mathbf{g}=N g(V, \Delta)$ stands for the total density of levels with given detuning and coupling and the contour $C$ has to be closed either in the upper or in the lower part of the complex plane $\Delta$ where the integrand vanishes.

In order to find an explicit expression for the averaged $\Psi$ we now perform the integration

$$
\Psi_{G}=\int e^{L_{G}} d \bar{\chi} d \chi d \bar{\chi}^{\prime} d \chi^{\prime}
$$

retaining in (8.106) only the terms $L_{G}$ of $L_{n}$ that contain Grassmann variables. Expanding $e^{L_{G}}$ in a Taylor series over each of the Grassmann variables up to first order and combining the terms proportional to $\bar{\chi} \chi \bar{\chi}^{\prime} \chi^{\prime}$, which according to the definition (7.177) is the only non-vanishing term remaining after the integration, we find

$$
\begin{align*}
\Psi_{G}= & \int d \bar{\chi} d \chi d \bar{\chi}^{\prime} d \chi^{\prime} \exp \left[-i \Delta_{\varepsilon} \bar{\chi} \chi+i \Delta_{\xi} \bar{\chi}^{\prime} \chi^{\prime}\right. \\
& \left.+\lambda \tau^{-1 / 2}\left(i \overline{\chi \chi}{ }^{\prime} w+i \bar{w} \chi^{\prime} \chi-x^{\prime} \bar{\chi} \alpha-\bar{\alpha} \chi \bar{x}^{\prime}-\bar{\chi}^{\prime} \bar{x} \beta-\bar{\beta} x \chi^{\prime}\right)\right] \\
= & \left(i \bar{x}^{\prime} x^{\prime} \alpha \bar{\alpha} \Delta_{\xi}-i \beta \bar{\beta} \bar{x} x \Delta_{\varepsilon}+\frac{i \lambda}{\sqrt{\tau}} w \bar{\alpha} \bar{\beta} x \bar{x}^{\prime}+\frac{i \lambda}{\sqrt{\tau}} x^{\prime} \bar{x} \beta \alpha \bar{w}\right) \frac{\lambda^{2}}{\tau} \\
& +\bar{x} x \bar{x}^{\prime} x^{\prime} \beta \bar{\beta} \alpha \bar{\alpha} \frac{\lambda^{4}}{\tau^{2}}+\Delta_{\xi} \Delta_{\varepsilon}-\frac{\lambda^{2}}{\tau} w \bar{w} . \tag{8.111}
\end{align*}
$$

The next step is the calculation of the integral

$$
\begin{equation*}
\Psi(\Delta)=\int e^{L^{\prime}} \Psi_{G} d \bar{x} d x d \bar{x}^{\prime} d x^{\prime} \tag{8.112}
\end{equation*}
$$

where the remaining terms

$$
\begin{aligned}
L^{\prime}= & -i \Delta_{\varepsilon} \bar{x} x+i \Delta_{\xi} \bar{x}^{\prime} x^{\prime}-\bar{x}_{0} \bar{x}_{0}^{\prime} x x^{\prime}+x_{0}^{\prime} \bar{x}_{0}^{\prime} \bar{x}^{\prime} x^{\prime} \\
& +x_{0} \bar{x}_{0} x \bar{x}-x_{0}^{\prime} x_{0} \bar{x}^{\prime} \bar{x}+\frac{i \lambda}{\sqrt{\tau}} \bar{x} x^{\prime} f+\frac{i \lambda}{\sqrt{\tau}} \bar{f} x \bar{x}^{\prime}
\end{aligned}
$$

of $L_{n}$ are taken into account with the allowance of (8.109) for the average over the couplings. With the help of the relations

$$
\begin{gathered}
\bar{x} x \bar{x}^{\prime} x^{\prime} \frac{\lambda^{2}}{\tau} e^{L^{\prime}}=-\frac{\partial^{2} e^{L^{\prime}}}{\partial \partial \bar{f}} ; i \bar{x}^{\prime} x^{\prime} e^{L^{\prime}}=\frac{\partial e^{L^{\prime}}}{\partial \Delta_{\xi}} ;-i \bar{x} x e^{L^{\prime}}=\frac{\partial e^{L^{\prime}}}{\partial \Delta_{\varepsilon}} \\
\frac{i \lambda}{\sqrt{\tau}} x \bar{x}^{\prime} e^{L^{\prime}}=\frac{\partial e^{L^{\prime}}}{\partial \bar{f}} ; \frac{i \lambda}{\sqrt{\tau}} x^{\prime} \bar{x} e^{L^{\prime}}=\frac{\partial e^{L^{\prime}}}{\partial f}
\end{gathered}
$$

one can write $\Psi(\Delta)$ in the form

$$
\begin{align*}
\Psi(\Delta)= & {\left[\left(\beta \bar{\beta} \frac{\Delta_{\varepsilon} \partial}{\partial \Delta_{\varepsilon}}+\alpha \bar{\alpha} \frac{\Delta_{\xi} \partial}{\partial \Delta_{\xi}}+\beta \alpha \bar{w} \frac{\partial}{\partial f}+w \bar{\alpha} \bar{\beta} \frac{\partial}{\partial \bar{f}}\right.\right.} \\
& \left.\left.-\frac{\partial^{2}}{\partial f \partial \bar{f}} \beta \bar{\beta} \alpha \bar{\alpha}-w \bar{w}\right) \frac{\lambda^{2}}{\tau}+\Delta_{\xi} \Delta_{\varepsilon}\right] \Xi, \tag{8.113}
\end{align*}
$$

where the Gaussian integral

$$
\Xi=\int e^{L^{\prime}} d \bar{x} d x d \bar{x}^{\prime} d x^{\prime}
$$

can be found explicitly by calculating the determinant of the quadratic form $\partial^{2} L^{\prime} / \partial x \partial x^{\prime}$, in agreement with (7.133). This yields

$$
\begin{equation*}
\Xi=\left[\left(i \Delta_{\xi} r-i \Delta_{\varepsilon} r^{\prime}-\Delta_{\varepsilon} \Delta_{\xi}-\frac{\lambda^{2}}{\tau} f \bar{f}\right)^{2}+4 r r^{\prime} \frac{\lambda^{2}}{\tau} f \bar{f}\right]^{-1 / 2} \tag{8.114}
\end{equation*}
$$

where the variables $r$ and $r^{\prime}$ stand for the combinations $x_{0} \bar{x}_{0}$ and $x_{0}^{\prime} \bar{x}_{0}^{\prime}$ respectively.

Now we perform the last remaining integration over the level position

$$
\begin{align*}
\int g \Psi d \Delta= & \int g d \Delta\left[\frac { \lambda ^ { 2 } } { \tau } \left(w \bar{\alpha} \bar{\beta} \frac{\partial}{\partial \bar{f}}+\beta \alpha \bar{w} \frac{\partial}{\partial f}-\beta \bar{\beta} \alpha \bar{\alpha} \frac{\partial^{2}}{\partial f \partial \bar{f}}\right.\right. \\
& \left.\left.-w \bar{w}+\beta \bar{\beta} \frac{\Delta_{\varepsilon} \partial}{\partial \Delta_{\varepsilon}}+\alpha \bar{\alpha} \frac{\Delta_{\xi} \partial}{\partial \Delta_{\xi}}\right)+\Delta_{\xi} \Delta_{\varepsilon}\right] \Xi \tag{8.115}
\end{align*}
$$

where $g=N \int g(V, \Delta) d^{2} V$ now denotes the uniform total density of the band levels, and find

$$
\begin{align*}
\int_{C} g \Psi d \Delta= & \int_{C} \Xi \Delta_{\xi} \Delta_{\varepsilon} d \Delta+\frac{\lambda^{2}}{\tau} \beta \bar{\beta} \int_{C} \frac{\Delta_{\varepsilon} \partial \Xi}{\partial \Delta_{\varepsilon}} d \Delta+\frac{\lambda^{2}}{\tau} \alpha \bar{\alpha} \int_{C} \frac{\Delta_{\xi} \partial \Xi}{\partial \Delta_{\xi}} d \Delta  \tag{8.116}\\
& +\frac{\lambda^{2}}{\tau}\left(\beta \bar{\beta}+\alpha \bar{\alpha}+w \bar{\alpha} \bar{\beta} \frac{\partial}{\partial \bar{f}}+\beta \alpha \bar{w} \frac{\partial}{\partial f}-\beta \bar{\beta} \alpha \bar{\alpha} \frac{\partial^{2}}{\partial f \partial \bar{f}}-w \bar{w}\right) \int_{C} g \Xi d \Delta .
\end{align*}
$$

The integral $\Lambda=\int_{C} g \Xi d \Delta$ can be given explicitly in terms of the hypergeometric function

$$
\begin{equation*}
\Lambda=\int_{C} \frac{g d y}{\sqrt{-\left(a-i y^{2}\right)^{2}-b^{2}}}=\sqrt{\frac{1}{a}} \pi g{ }_{2} F_{1}\left(\frac{1}{4}, \frac{3}{4} ; 1, \frac{b^{2}}{a^{2}}\right) \tag{8.117}
\end{equation*}
$$

with

$$
\begin{gather*}
a=\frac{\left(r-r^{\prime}\right)^{2}}{4}+\frac{\lambda^{2}}{\tau} f \bar{f}+i \frac{r+r^{\prime}}{2}\left(\Delta_{\varepsilon}-\Delta_{\xi}\right)-\frac{1}{4}\left(\Delta_{\xi}-\Delta_{\varepsilon}\right)^{2} \\
b=\sqrt{\frac{4 r r^{\prime} \lambda^{2} f \bar{f}}{-\tau}} \tag{8.118}
\end{gather*}
$$

whereas the remaining integrals

$$
\begin{gather*}
\Lambda_{\varepsilon}=\int \frac{\Delta_{\varepsilon} \partial}{\partial \Delta^{\prime}} \Xi d \Delta=\frac{r^{\prime}-r}{2 i} \frac{\partial \Lambda}{\partial X}+\frac{X}{2} \frac{\partial \Lambda}{\partial X}-\frac{1}{2} \Lambda \\
\Lambda_{\xi}=\int \frac{\Delta_{\xi} \delta}{\partial \Delta_{\xi}} \Xi d \Delta=\frac{r-r^{\prime}}{2 i} \frac{\partial \Lambda}{\partial X}+\frac{X}{2} \frac{\partial \Lambda}{\partial X}-\frac{1}{2} \Lambda  \tag{8.119}\\
\quad \int d \Delta \Xi \Delta_{\xi} \Delta_{\varepsilon}=\Phi-\frac{\left(r-r^{\prime}\right)^{2}}{4} \Lambda-\frac{X^{2}}{4} \Lambda
\end{gather*}
$$

can also be expressed in term of the hypergeometric functions

$$
\begin{equation*}
\Phi=\int_{C} \frac{g y^{2} d y}{\sqrt{-\left(a-i y^{2}\right)^{2}-b^{2}}}=-\sqrt{a} \pi g{ }_{2} F_{1}\left(-\frac{1}{4}, \frac{1}{4} ; 1, \frac{b^{2}}{a^{2}}\right) \tag{8.120}
\end{equation*}
$$

where $X=\Delta_{\varepsilon}-\Delta_{\xi}=F(\varepsilon)-F(\xi)+\xi-\varepsilon$. Note that the hypergeometric functions have branching points at $b^{2} / a^{2}=1$, and the proper choice of the branches of $\Lambda$ and $\Phi$ depend on the sign of the imaginary part of $a-b$, that is on the sign of $\zeta$. Equation (8.116) finally adopts the form

$$
\begin{align*}
\int g \Psi d \Delta= & \beta \bar{\beta} \frac{\lambda^{2}}{\tau} \Lambda_{\xi}+\alpha \bar{\alpha} \frac{\lambda^{2}}{\tau} \Lambda_{\varepsilon}-\frac{\lambda^{2}}{\tau} \frac{\partial^{2} \Lambda}{\partial f \partial \bar{f}} \beta \bar{\beta} \alpha \bar{\alpha}+w \bar{\alpha} \bar{\beta} \frac{\lambda^{2} \partial \Lambda}{\tau \partial \bar{f}} \\
& +\beta \alpha \bar{w} \frac{\lambda^{2} \partial \Lambda}{\tau \partial f}-w \bar{w} \frac{\lambda^{2} \Lambda}{\tau}+\Phi-\frac{\left(r-r^{\prime}\right)^{2}}{4} \Lambda-\frac{X^{2}}{4} \Lambda . \tag{8.121}
\end{align*}
$$

We now substitute (8.121) into (8.111) and perform the exact integrations over the variables $d \bar{\alpha} d \alpha d \bar{\beta} d \beta$ starting with the Grassmann variables. Expanding in a Taylor series and retaining only the terms linear in $\bar{\alpha}, \alpha, \bar{\beta}, \beta$ we arrive at

$$
\begin{align*}
\left\langle\rho_{0}(\varepsilon, \xi)\right\rangle= & \int \frac{e^{-\tau}}{\tau} d \tau d w d \bar{w} d f d \bar{f} \int d r d r^{\prime} \\
& {\left[\left(1-\frac{\lambda^{2}}{\tau} \Lambda_{\xi}\right)\left(1-\frac{\lambda^{2}}{\tau} \Lambda_{\varepsilon}\right)-\frac{\lambda^{2}}{\tau} \frac{\partial^{2} \Lambda}{\partial f \partial \bar{f}}+\bar{w} w \frac{\partial \Lambda \lambda^{2}}{\tau \partial \bar{f}} \frac{\partial \Lambda \lambda^{2}}{\tau \partial f}\right] }  \tag{8.122}\\
& \exp \left(i \varepsilon r-i \xi r^{\prime}-w \bar{w} \frac{\Lambda \lambda^{2}}{\tau}-w \bar{w}-f \bar{f}+\Phi-\frac{\left(r-r^{\prime}\right)^{2}}{4} \Lambda-\frac{X^{2}}{4} \Lambda\right),
\end{align*}
$$

where we have taken into account that

$$
\frac{1}{\pi^{2}} \int_{-\infty}^{\infty} \ldots d \bar{x}_{0} d x_{0} d \bar{x}_{0}^{\prime} d x_{0}^{\prime} \rightarrow \int_{0}^{\infty} \ldots d r d r^{\prime}
$$

After integration over $d w d \bar{w}$ this yields the final general and exact expression

$$
\begin{align*}
\left\langle\rho_{0}(\varepsilon, \xi)\right\rangle= & \int \frac{e^{-\tau}}{\tau} d \tau d f d \bar{f} \int d r d r^{\prime}\left[\frac{-\tau}{\Lambda \lambda^{2}+\tau} \frac{\lambda^{2}}{\tau} \frac{\partial^{2} \Lambda}{\partial f \partial \bar{f}}\right. \\
& \left.+\frac{\tau}{\Lambda \lambda^{2}+\tau}\left(1-\frac{\lambda^{2}}{\tau} \Lambda_{\xi}\right)\left(1-\frac{\lambda^{2}}{\tau} \Lambda_{\varepsilon}\right)+\left(\frac{\tau}{\Lambda \lambda^{2}+\tau}\right)^{2} \frac{\partial \Lambda \lambda^{2}}{\tau \partial \bar{f}} \frac{\partial \Lambda \lambda^{2}}{\tau \partial f}\right] \\
& \exp \left(i \varepsilon r-i \xi r^{\prime}-f \bar{f}+\Phi-\frac{\left(r-r^{\prime}\right)^{2}}{4} \Lambda-\frac{X^{2}}{4} \Lambda\right) \tag{8.123}
\end{align*}
$$

for the ensemble averaged population of the state $|0\rangle$. Further analysis should rely on the particular choice of the functions $F(\varepsilon)$ and $\lambda^{2}(\varepsilon, \xi)$ that depend on the particular type of random walk.

### 8.2.3 Three Types of Random Walk and the Asymptotic Decay

As we have already mentioned, the analysis of population dynamics in the long-time limit, resulting from the quantum interference of different selfintersecting trajectories, depends on particular properties of the random
walk. The quantities $G(\varepsilon)$ and $\lambda^{2}(\varepsilon, \xi)$ responsible for the specific behavior of the random walk and entering (8.121) and (8.123) have a clear physical meaning for the system of band levels $|n\rangle$ and the states $|\kappa\rangle$ of the continuum that does not include the state $|0\rangle$. The quantity $G(\varepsilon)$ has been already introduced in Sect.3.2.1 for the level-band problem, and it allows for the level probability amplitude decay rate to the continuum. Now we consider $G(\varepsilon)$ as a characteristic of the band levels interacting with the continuum responsible for the random walk. The quantity $\lambda^{2}(\varepsilon, \xi)$ represents the Fourier transform of the time correlation of the probability amplitudes to be in a given state of the band under the condition that at the initial time moment $t=0$, the population was localized in another level of the band. After being integrated over the variable $\eta=(\varepsilon+\xi) / 2$, this quantity depends only on the difference $\zeta=\varepsilon-\xi$ and represents the Fourier transform of the probability to be in the given state.

We consider two rather general particular cases of returnable and nonreturnable random walks. For a returnable random walk, the Fourier transform $\lambda^{2}(\varepsilon, \xi)$ diverges at the limit $\varepsilon \rightarrow \xi$ according to the power law $(\varepsilon-\xi)^{-\nu}$ where the power index $\nu$ also describes the increase of the expectation number of returns $\mathcal{N} \sim t^{\nu}$ at time $t$. For a non-returnable random walk, $\lambda \rightarrow 0$ as $\varepsilon \rightarrow \xi$. For the first case we also identify two different particular cases: a random walk with a finite correlation time $\tau_{r}$ and a $\delta$-correlated random walk. The correlation time of the random walk is a typical residence time $\tau_{r}$, during which the probability amplitude of the quantum particle remains at a state of the band, once it has returned to the state from the continuum. This time is associated with dephasing of the probability amplitudes, whereas the absolute value of the amplitude changes smoothly. The width of the quantity $\lambda^{2}(\zeta, \eta)$ as a function of $\eta$ is inversely proportional to $\tau_{r}$. The order of magnitude of $\tau_{r}$ can be found from the second derivative

$$
\begin{equation*}
\frac{\partial^{2} \lambda^{2}(\eta)}{2 \partial \eta^{2}}=-\frac{\lambda^{2}(\eta)}{\tau_{r}^{2}} \tag{8.124}
\end{equation*}
$$

at the point $\eta=0$. It also relates to the width of the continuum seen by each state of the band. The function $\lambda^{2}(\zeta, \eta)$ does not depend on $\eta$ at all for a $\delta$-correlated random walk, when the particle wavepacket rapidly comes and leaves the band state, keeping the mean population unchanged during some short time interval.

## Non-Returnable Random Walk

For $\lambda \rightarrow 0$, the only important singularity of the integrand (8.123) as a function of $\tau$ is the simple pole at $\tau=0$, and for this case the integration over $d \tau$ yields

$$
\begin{gather*}
\left\langle\rho_{0}(\varepsilon, \xi)\right\rangle=\int d f d \bar{f} \int d r d r^{\prime} \\
\exp \left(i \varepsilon r-i \xi r^{\prime}-f \bar{f}-\sqrt{a} \pi g-\frac{\left(r-r^{\prime}\right)^{2}}{4 \sqrt{a}} \pi g-\frac{X^{2}}{4 \sqrt{a}} \pi g\right) . \tag{8.125}
\end{gather*}
$$

with

$$
a=\frac{\left(r-r^{\prime}\right)^{2}}{4}+i \frac{r+r^{\prime}}{2} X-\frac{X^{2}}{4}
$$

where we have employed the fact that hypergeometric functions tend to one at zero argument. For the long-time limit, when $\varepsilon \sim \xi \sim 1 / t$ the parameter $X=F(\varepsilon)-F(\xi)+\xi-\varepsilon \simeq-2 i \gamma_{c}$ can be considered as independent of $\varepsilon$ and $\xi$, and it is given by the decay rate $\gamma_{c}$ of the amplitudes of the band states decaying to the continuum. In this case, the inverse Fourier transformation over $\varepsilon$ and $\xi$ immediately gives the Dirac $\delta$-functions $\delta(r-t)$ and $\delta\left(r^{\prime}-t\right)$, that is $a=2 \gamma t$, and we arrive at

$$
\begin{equation*}
\left\langle\rho_{0}(t)\right\rangle=\int d f d \bar{f} \exp \left(-f \bar{f}-\sqrt{a} \pi g-\frac{\gamma^{2}}{4 \sqrt{a}} \pi g\right) \simeq \exp \left(-\pi g \sqrt{\gamma_{c} t}\right) \tag{8.126}
\end{equation*}
$$

or $\left\langle\rho_{0}(t)\right\rangle \simeq \exp \left(-\pi g \sqrt{\left\langle V^{2}\right\rangle \gamma_{c} t}\right)$ for dimensional variables.
The exponent of the square root time dependence emerges from the ensemble average over the randomly distributed energies of the levels in the band. Namely, it results from the unlikely long leaving realizations of the system when the level $|0\rangle$ is separated from the closest level $|n\rangle$ of the band by a big energy gap $\Delta$, as shown in Fig $.8 .9(\mathrm{~b})$. The probability of such a realization for the independently distributed $\Delta_{n}$ is exponentially small $\sim \exp (-g \Delta)$. However the contribution of this realization is dominating for large $t$ since it is long-leaving, with a decay rate inversely proportional to $\Delta$. Indeed, the composite matrix element $V_{0 n} V_{n \kappa} / \Delta$ of the transition from the level $|0\rangle$ to a state $|\kappa\rangle$ of the continuum via a level $|n\rangle$ of the band detuned from the state $|0\rangle$ by the energy difference $\Delta_{n}$, gives the contribution $\left(V_{0 n} V_{n \kappa} / \Delta\right)^{2} g_{\kappa}=\left(V_{0 n} / \Delta_{n}\right)^{2} \gamma_{c}$ to the decay rate. The contribution of all levels with the detuning larger than a gap size $\Delta$ is given by the integral $\int_{\Delta}^{\infty}\left(V_{0 n} / \Delta_{n}\right)^{2} \gamma_{c} g d \Delta_{n}$ that is a quantity $\gamma_{c} g V^{2} / \Delta$ which is small for a large gap size and which results in a decay $\exp \left(-t \gamma_{c} g V^{2} / \Delta\right)$ of $\rho_{0}(t)$. Averaging over the gap size yields

$$
\begin{equation*}
\left\langle\rho_{0}(t)\right\rangle=\int d \Delta \exp \left(-t \gamma_{c} g V^{2} / \Delta-g \Delta\right) \sim \exp \left(-\pi g \sqrt{V^{2} \gamma_{c} t}\right) \tag{8.127}
\end{equation*}
$$

where we have ignored a factor in front of the exponent, slowly dependent on time. In the units where $V=1$ this expression coinsides with (8.126).

## Returnable Random Walk

Now let us consider the technically more difficult case of a returnable random walk. For $\lambda \neq 0$ the point $\tau=0$ gives no contribution to the integral (8.123), since the hypergeometric function entering the term $\Phi$ in the exponent tends at this point to infinity. The contribution comes now from the singularity at the point $\tau=-\Lambda \lambda^{2}$. In order to avoid cumbersome calculations of the
derivatives of hypergeometric functions associated with the second-order pole we make the replacement of the integration variables $f \rightarrow-\tau f$ which moves the variable $\tau$ out of the arguments of $\Lambda$ and $\Phi$. Such a replacement does not affect the convergence of the integrals over $d f d f$, provided the term $\Lambda \lambda^{2}$ has a positive real part. If this is not the case, the replacement should be taken with the opposite sign. The integral (8.123) takes the form

$$
\begin{align*}
\left\langle\rho_{0}(\varepsilon, \xi)\right\rangle= & \int e^{-\tau} d \tau d f d \bar{f} \int d r d r^{\prime}\left[\frac{1}{\Lambda \lambda^{2}+\tau} \frac{\partial^{2} \Lambda \lambda^{2}}{\tau \partial \partial \partial \bar{f}}\right. \\
& \left.-\frac{\tau}{\Lambda \lambda^{2}+\tau}\left(1-\frac{\lambda^{2}}{\tau} \Lambda_{\xi}\right)\left(1-\frac{\lambda^{2}}{\tau} \Lambda_{\varepsilon}\right)-\left(\frac{1}{\Lambda \lambda^{2}+\tau}\right)^{2} \frac{\partial \Lambda \lambda^{2}}{\partial \bar{f}} \frac{\partial \Lambda \lambda^{2}}{\tau \partial f}\right] \\
& \exp \left(i \varepsilon r-i \xi r^{\prime}+\tau f \bar{f}+\Phi-\frac{\left(r-r^{\prime}\right)^{2}}{4} \Lambda-\frac{X^{2}}{4} \Lambda\right) \tag{8.128}
\end{align*}
$$

where the arguments of $\Phi$ and $\Lambda$ read

$$
\begin{gather*}
a=\frac{\left(r-r^{\prime}\right)^{2}}{4}-\lambda^{2} f \bar{f}+i \frac{r+r^{\prime}}{2} X-\frac{X^{2}}{4}  \tag{8.129}\\
b=\sqrt{4 r r^{\prime} \lambda^{2} f \bar{f}}
\end{gather*}
$$

The integration over $d \tau$ yields

$$
\begin{align*}
\left\langle\rho_{0}(\varepsilon, \xi)\right\rangle= & \int d f d \bar{f} \int d r d r^{\prime}\left[-\frac{\partial^{2} \Lambda}{\Lambda \partial f \partial \bar{f}}+\Lambda \lambda^{2}\left(1+\frac{\Lambda_{\xi}}{\Lambda}\right)\left(1+\frac{\Lambda_{\varepsilon}}{\Lambda}\right)\right. \\
& \left.+\frac{\partial \Lambda}{\Lambda \partial \bar{f}} \frac{\partial \Lambda}{\Lambda \partial f}\left(1-\Lambda \lambda^{2}+f \bar{f} \Lambda \lambda^{2}\right)\right]  \tag{8.130}\\
& \times \exp \left(i \varepsilon r-i \xi r^{\prime}+\Phi-\frac{\left(r-r^{\prime}\right)^{2}}{4} \Lambda-\frac{X^{2}}{4} \Lambda+\Lambda \lambda^{2}-\Lambda \lambda^{2} f \bar{f}\right) .
\end{align*}
$$

No further significant simplifications can be performed in the general form of (8.130), even by introducing a single integration variable $u$ instead of the product $f \bar{f}$ of two variables. However one can pursue the simplification by making the replacement $f \bar{f} \rightarrow u$ in the long-time limit, where the return probability $\lambda^{2}$ is asymptotically large and hence the asymptotic form

$$
\begin{gather*}
\Lambda\left(r, r^{\prime}, u, \lambda^{2}\right) \rightarrow \frac{\Lambda}{2 u^{1 / 4} \lambda^{1 / 2} \sqrt{r r^{\prime}}} ; \Lambda=\pi g \sqrt{i \pi} \Gamma\left(\frac{1}{4}\right)^{-2}  \tag{8.131}\\
\Phi\left(r, r^{\prime}, u, \lambda^{2}\right) \rightarrow-2 \Phi u^{1 / 4} \lambda^{1 / 2} \sqrt{r r^{\prime}} ; \Phi=\pi g \sqrt{i^{3} \pi} \Gamma\left(\frac{3}{4}\right)^{-2}
\end{gather*}
$$

can be employed for the functions $\Phi\left(r, r^{\prime}, u, \lambda^{2}\right)$ and $\Lambda\left(r, r^{\prime}, u, \lambda^{2}\right)$. Note that the functions are now independent of $X$. Hereafter the symbols $\Phi$ and $\Lambda$ without arguments denote the large constants given by (8.131) in terms of $\Gamma$-functions. Substitution of (8.131) into (8.130) with allowance of (8.119) yields

$$
\begin{align*}
\left\langle\rho_{0}(\varepsilon, \xi)\right\rangle= & \int d u \int d r d r^{\prime} \frac{\Lambda \lambda^{3 / 2}(3+u)}{32 u^{5 / 4} \sqrt{r r^{\prime}}} e^{i \varepsilon r-i \xi r^{\prime}} \\
& \exp \left[-2 \Phi u^{1 / 4} \lambda^{1 / 2} \sqrt{r r^{\prime}}-\Lambda \frac{\left(r-r^{\prime}\right)^{2}+X^{2}-4 \lambda^{2}+4 \lambda^{2} u}{8 u^{1 / 4} \lambda^{1 / 2} \sqrt{r r^{\prime}}}\right] . \tag{8.132}
\end{align*}
$$

We neglect the finite value $X^{2}$ compared to the large function $\lambda^{2}$ and employ the hyperbolic coordinates $R, \theta, \phi$ introduced by the relations

$$
\begin{equation*}
r=R e^{\theta+\phi} \frac{\Lambda}{2 \Phi} ; r^{\prime}=R e^{-\theta-\phi} \frac{\Lambda}{2 \Phi} ; u=\frac{e^{4 \phi}}{\lambda^{2}}\left(\frac{\Lambda}{2 \Phi}\right)^{2} \tag{8.133}
\end{equation*}
$$

with the Jacobian $-e^{6 \phi} R \Lambda^{4} / 2 \lambda^{2} \Phi^{4}$. This results in

$$
\begin{align*}
& \left\langle\rho_{0}(\zeta, \eta)\right\rangle=\int d \phi\left(e^{4 \phi}+12 \bar{\lambda}^{2}\right) \int d R \exp \left[\left(-\sqrt{2} e^{\phi} R-\frac{4 e^{-2 \phi} \bar{\lambda}^{2}-e^{2 \phi}}{\sqrt{2} R}\right) \frac{\bar{g} U}{2}\right] \\
& \frac{i(\bar{g} U)^{3}}{16} \int d \theta \exp \left[\left(i e^{\phi} 2 \eta \sinh [\theta]-\sqrt{2} \zeta e^{\phi} \cosh [\theta]-\sinh [\theta]^{2}\right) \frac{R \bar{g} U}{2}\right], \tag{8.134}
\end{align*}
$$

where $\bar{g}=\sqrt{\Lambda \Phi}=\sqrt{\pi} g / \sqrt{2}, U=\Phi \Lambda^{-1}, \lambda=2 \bar{\lambda} U, \zeta=i \sqrt{2} \bar{g}(\varepsilon-\xi)$, and $\eta=\bar{g}(\varepsilon+\xi) / 2$. By the last two substitutions we have scaled the time variable $t \rightarrow t / \bar{g}$, which means that time is now measured in the band recurrence periods (Heisenberg times) discussed in Sect.3.3.1. We have also rotated the variable $\varepsilon-\xi$ in the complex plane such that the integration contour of the inverse Fourier transformation will circumvent the positive part of the real axis. The Fourier transformation over this variable thus becomes the Laplace transformation. The expression is multiplied by the Jacobian $i \sqrt{2} \bar{g}^{2}$ of the last two replacements as well.

Since the parameter $\bar{g}$ is large, the integration over $d \theta$ in (8.134) can be performed by the saddle point method. Expanding the function in the exponent near the point $\theta=0$, one finds the first integral over the hyperbolic angle $\theta$

$$
\begin{equation*}
\frac{\sqrt{2}}{\sqrt{R}}\left(\frac{\pi \sqrt{2}}{U \bar{g}}\right)^{1 / 2}\left(1+i e^{\phi} \zeta\right)^{-1 / 2} \exp \left[\frac{-R U}{\sqrt{2}} \bar{g}\left(e^{\phi} \zeta+\frac{e^{2 \phi} \eta^{2}}{1+\zeta e^{\phi}}\right)\right] . \tag{8.135}
\end{equation*}
$$

The second integral over the hyperbolic radius $R$ now has the structure

$$
\int_{0}^{\infty} \frac{d R}{\sqrt{R}} \exp \left(-A R-\frac{1}{R} B\right)=\sqrt{\frac{\pi}{A}} e^{-2 \sqrt{A B}}
$$

and yields

$$
\begin{align*}
\left\langle\rho_{0}(\zeta, \eta)\right\rangle= & \frac{2^{1 / 4} \pi i U^{2}}{8} \int d \phi \frac{e^{-\phi}\left(e^{4 \phi}+12 \bar{\lambda}^{2}\right)}{\sqrt{\eta^{2}+\zeta^{2}+1+2 \zeta \cosh \phi}} \\
& \exp \left[-\frac{U \bar{g}\left(4 \bar{\lambda}^{2}-e^{4 \phi}\right)^{1 / 2}}{\left(1+\zeta e^{\phi}\right)^{1 / 2}}\left(\eta^{2}+\zeta^{2}+1+2 \zeta \cosh \phi\right)^{1 / 2}\right] \tag{8.136}
\end{align*}
$$

One can neglect the small quantity $\zeta$ everywhere in (8.136) except in those terms where it is multiplied by exponentially large factors. We also note that
for small $\zeta$, two different contributions to the integral come from two distinct domains of large positive and large negative $\phi$, where only one of the two exponents entering $\cosh \phi$ can be retained. This suggests the substitution $x=\zeta e^{\phi}$ for large positive $\phi$ and the substitution $x=\zeta e^{-\phi}$ for large negative $\phi$ with $x$ ranging from 0 to $\infty$, such that the integral adopts the form

$$
\begin{align*}
\left\langle\rho_{0}(\zeta, \eta)\right\rangle= & \frac{2^{1 / 4} \pi i U^{2}}{8} \int_{0}^{\infty} d x \frac{\left(x / \zeta^{2}+12 \bar{\lambda}^{2} \zeta^{2} / x^{2}\right)}{\zeta \sqrt{\eta^{2}+1+x}} \\
& \times \exp \left[-\frac{U \bar{g}\left(4 \bar{\lambda}^{2}-\zeta^{-4} x\right)^{1 / 2}}{(1+x)^{1 / 2}}\left(\eta^{2}+1+x\right)^{1 / 2}\right]  \tag{8.137}\\
& +\int_{0}^{\infty} d x \frac{3 \bar{\lambda}^{2} 2^{1 / 4} \pi i U^{2}}{2 \zeta \sqrt{\eta^{2}+1+x}} \cdot \exp \left[-2 U \bar{g}\left(\bar{\lambda}^{2}\right)^{1 / 2}\left(\eta^{2}+1+x\right)^{1 / 2}\right] .
\end{align*}
$$

Inspection shows that the first integral contributes only at short times and can be ignored. The second integral can be evaluated exactly and yields the probability to be in the state $|0\rangle$ for the returnable random walk

$$
\begin{equation*}
\left\langle\rho_{0}(\zeta, \eta)\right\rangle=\frac{3 \sqrt{\pi} \lambda(\zeta, \eta)}{2^{1 / 4} \zeta g V} \exp \left[-\sqrt{2 \pi}\left(\eta^{2}+V^{4} g^{2}\right)^{1 / 2} \lambda(\zeta, \eta)\right] . \tag{8.138}
\end{equation*}
$$

In this final result for the general case of the returnable random walk, we have restored the original notation $\lambda^{2}(\zeta, \eta)=4 \bar{\lambda}^{2} U^{2}$ for the Fourier transform of the random walk correlation function and the density of the band states $\bar{g}=\sqrt{\pi} g / \sqrt{2}$. The mean squared coupling $V=\left\langle V^{2}\right\rangle^{1 / 2}$ is also included.

## The Role of the Correlation Time for a Returnable Walk

We now analyze the role of the correlation time $\tau_{r}$. For a $\delta$-correlated random walk, the Fourier transform $\lambda^{2}(\zeta, \eta)$ does not depend on $\eta$. Therefore the main contribution to the inverse Fourier transformation over $d \eta$ comes from the domain $\eta=0$ where the exponent in (8.138) has the only saddle point. Evaluation of the integral results in

$$
\begin{equation*}
\left\langle\rho_{0}(\zeta)\right\rangle=\frac{3 \pi^{3 / 4} \sqrt{\lambda(\zeta)}}{\zeta \sqrt{g}} \exp \left[-\sqrt{2 \pi} V^{2} g \lambda(\zeta)\right] \tag{8.139}
\end{equation*}
$$

and for a random walk with the power law dependence $\lambda(\zeta)=\kappa \zeta^{-\alpha}$ at the limit of small $\zeta$ this gives the time dependence

$$
\begin{align*}
\left\langle\rho_{0}(t)\right\rangle & =\int \frac{3 \pi^{3 / 4} \sqrt{\lambda(\zeta)}}{2 \pi \zeta \sqrt{g}} e^{-\zeta t-\sqrt{2 \pi} V^{2} g \lambda(\zeta)} d \zeta \\
& \sim \frac{t^{1 /(1+\alpha)}}{g^{(3+2 \alpha) /(1+\alpha)}} \exp \left[- \text { const } t^{\alpha /(1+\alpha)}\right] \tag{8.140}
\end{align*}
$$

For the case of a finite $\tau_{r}$ we employ the model $\lambda^{2}(\zeta, \eta)=\lambda^{2}(\zeta) /\left(\eta^{2}+\gamma^{2}\right)$ where $\gamma=1 / \tau_{r}$ determines the width associated with the dephasing. After the replacement

$$
\begin{equation*}
\sqrt{2 \pi} \sqrt{\frac{\eta^{2}+V^{4} g^{2}}{\left(\eta^{2}+\gamma^{2}\right)}} \lambda(\zeta) \rightarrow Z \tag{8.141}
\end{equation*}
$$

equation (8.138) yields

$$
\begin{equation*}
\left\langle\rho_{0}(\zeta)\right\rangle=\int_{0}^{\infty} \frac{-i 32^{1 / 4} \pi \lambda^{2}(\zeta) Z \sqrt{\gamma^{2}-V^{4} g^{2}} e^{-Z}}{\zeta g V\left(Z^{2}+2 \pi \lambda^{2}(\zeta)\right) \sqrt{Z^{2} \gamma^{2}+2 \pi \lambda^{2}(\zeta) V^{4} g^{2}}} d Z \tag{8.142}
\end{equation*}
$$

Here we have employed the fact that the integration contour for $\eta$ can be closed in the upper part of the complex plane (which was impossible for the $\delta$-correlating random walk), whereas the lower and upper integration limits for $Z$ correspond to the points $\eta=i V^{2} g$ and $\eta=i \gamma$, respectively. The analytical property of $\lambda^{2}(\zeta, \eta)$ that makes such a transformation of the contour possible has an important consequence - contrary to (8.139) the integral (8.142) contains a term which corresponds to a population, nonvanishing in the long-time limit. Indeed, for small values of the parameter $V^{4} g^{2} \gamma^{-2}$, by neglecting $Z^{2}$ relative to the large function $2 \pi \lambda^{2}(\zeta)$ we find the integral exactly in terms of the Bessel functions which we however give only in the asymptotic form for large $\lambda^{2}(\zeta)$

$$
\begin{align*}
\left\langle\rho_{0}(\zeta)\right\rangle= & \frac{32^{-3 / 4}}{i \zeta g V} \int_{0}^{\infty} \frac{Z e^{-Z}}{\sqrt{Z^{2}+2 \pi \lambda^{2}(\zeta) V^{4} g^{2} \gamma^{-2}}} d Z  \tag{8.143}\\
& \frac{32^{3 / 4}}{i \zeta g V}+\frac{3 i}{\zeta} \sqrt{\frac{\lambda(\zeta)}{\sqrt{\pi} \gamma g}} \sin \left(\frac{\lambda(\zeta) V^{2} g}{\gamma}\right) \ln \left(\frac{\sqrt{2 \pi} \lambda(\zeta) V^{2} g}{\gamma}\right)
\end{align*}
$$

in order to be consistent within the approximation.
The asymptotic regime corresponding to (8.143) is attained only for very long times, when the number of returns $\sim \lambda^{2}(1 / t)$ exceeds the number of levels $(\gamma g)^{2}$ in the spectral band of width $\gamma=1 / \tau_{r}$ corresponding to the correlation time. Therefore it is never attained for the $\delta$-correlated random walk with $\tau_{r}=0$.

The inverse Laplace transformation of (8.143) yields

$$
\begin{equation*}
\left\langle\rho_{0}(t)\right\rangle=\frac{32^{3 / 4}}{g V}+\frac{\text { const } t^{1 /(1+\alpha)}}{(g / \gamma)^{(3+2 \alpha) /(1+\alpha)}} \exp \left[- \text { const }(g / \gamma)^{-1 /(1+\alpha)} t^{\alpha /(1+\alpha)}\right] \tag{8.144}
\end{equation*}
$$

where the last term is estimated for $\lambda(\zeta)=\kappa \zeta^{-\alpha}$ by analogy to (8.140). One sees, that the asymptotic population is given by the inverse of the expectation value of resonances $g V$ for the level-band system and does not contain any
trace of the correlation time, which only determines the typical time when the constant asymptotic population of the state $|0\rangle$ is reached.

Summarizing the results obtained, we can say that a level does not decay completely to a band corresponding to a returnable random walk with a finite time of the correlation of the probability amplitudes. It keeps a residual population given by the inverse expectation number of resonances $g V$. For a non-returnable random walk, the decay manifests a universal time dependence $\exp (-\sqrt{t})$, whereas for a $\delta$-correlated returnable random walk, the decay law $\exp \left[-t^{\alpha /(1+\alpha)}\right]$ is governed by the power index $\alpha$ of the Fourier transformed return probability $\lambda(\zeta)=\kappa \zeta^{-\alpha}$.

### 8.3 Manifestation of Quantum Complexity in the State Density

Thus far we were mainly focused on the dynamics of populations in complex quantum systems and the time evolution of probability amplitudes was considered only as an auxiliary tool in this context. However, some important physical properties such as the density of states and the linear susceptibility of quantum objects at a given frequency are primarily related to the dynamics of the probability amplitudes.

Generally speaking, all of the physical properties that can be expressed as a linear functional of the Green's operator of a system, are associated with the dynamical and spectral properties of the probability amplitudes. Moreover, all results obtained for the ensemble averaged amplitudes $\langle\psi(\varepsilon)\rangle$ corresponding to the initial condition $\psi(t=0)=1$ are valid for the Green's functions $G_{0 n}(E)=\psi_{n}(E)$, since in fact these two objects are identical. In this section, we consider a general rule which governs the transformation of the spectral density and other physical characteristics of a large number of interacting quantum states, linear in $G(E)$. The Fourier transforms of such spectral characteristics give the corresponding time-dependent response functions. The universal features in these spectral dependences, typical of complex multilevel quantum systems, give rise to some universality in the responses, that at certain conditions can manifest a universal behavior of physically different response functions at asymptotically long time.

We start by addressing the question: what happens to the density of states of an arbitrary multilevel quantum system when we couple the levels by random interaction? The rigorous analysis of this problem relies on involved calculations, but the final result which governs the reconstruction of the spectra in the limit of a large number of quantum states mixed up by a random perturbation has a very simple, although non-trivial form. It is closely related with the procedure of renormalization that has been employed in Sect.4.3.1, in the context of the analysis of population dynamics for time intervals shorter than the inverse density of states. We can also illustrate this
result by simple heuristic speculation, based on second-order perturbation theory. Several examples of the transformation of the state density and the absorption line profiles follow.

The rigorous analysis, for the case when the interaction mixes up just a small number of states $g V \sim 1$, is a more challenging task. One has to take into account the role of recurrences and revivals, which manifests itself in the dominating contribution of the trajectories with multiple self-intersections, similar to the trajectories considered in the previous section in the context of the incomplete level decay, and population localization. In this section we give a description based directly on the physically more meaningful calculation of the self-intersecting trajectories and the corresponding diagrams, instead of the more compact approach of Sect.8.2 involving the Grassmann integration. We try to show the parallels between these two approaches when it is possible.

### 8.3.1 Spectrum Transformation Induced by Random Perturbation

In the second-order perturbation theory, a weak coupling $\hat{V}$ shifts the energies $\Delta_{l}$ of levels such that they acquire the new values

$$
\begin{equation*}
\tilde{\Delta}_{l}=\Delta_{l}+\sum_{m} \frac{\left|V_{l, m}\right|^{2}}{\Delta_{l}-\Delta_{m}} \tag{8.145}
\end{equation*}
$$

This implies that the physical properties associated with a given level undergo a shift to a different spectral region. Let us take, for example, the linear susceptibility $f_{l}(E)=\mu_{l}^{2}\left(E-\Delta_{l}-i \gamma_{l}\right)^{-1}$ corresponding to the level $l$. Here $\gamma_{l}$ is the decay rate and $\mu_{l}$ the dipole moment of the transition to this level from the ground state. The real part of the susceptibility $\operatorname{Re} f_{l}\left(E ; \Delta_{l}\right)$ gives the Raman scattering amplitude, and the imaginary part $\operatorname{Im} f_{l}\left(E ; \Delta_{l}\right)$ is the absorption cross-section. For the perturbed spectra in these expressions we have to replace $\Delta_{l}$ by $\tilde{\Delta}_{l}$.

But what happens when the interaction is not small? Evidently, the level shifts become larger. Besides the shifts, the strong interaction mixes up the quantum states, so that the "individual properties" of each of the levels become distributed among a number of energy eigenstates. It turns out that for a random perturbation, this spreading of the "individual properties" does not manifest itself in the resulting spectra: after the ensemble average, only the level shifts remain important. These shifts may however contain imaginary components, thus moving the state responsible for a certain contribution to the quantity of interest to the complex plane of energies. In other words, a physical value $\Phi(E)$, which had a given magnitude at the energy $E$, now has the same magnitude at the energy $\tilde{E}$. Here the function $\Phi(E)$ is supposed to be proportional to the Green's function, that is, it corresponds to a linear property of the system.

## A Simple Rule for Large Complex Systems

The main features of the transformation of the density of states follows from a procedure almost identical to the renormalization discussed in Sect.4.1.3. It only has to be generalized to an arbitrary density of the unperturbed states. This gives the relation between $E$ and $\tilde{E}$, that is between the domain of the unperturbed energy spectrum associated with the contribution to the particular physical property, and the new energy position of the state responsible for this contribution. We note that as a result of such a transformation, the new energy position could, and usually does, acquire an imaginary part. In other words, the quantum states contributing in a specific way to the physical property of the system, are not necessarily energy eigenstates. In the limit of a large number of levels, this energy is given by the transcendental algebraic equation resembling (8.145), which reads

$$
\begin{equation*}
E=\tilde{E}+\sum_{m} \frac{\left\langle V^{2}\right\rangle}{\tilde{E}-\Delta_{m}-i \gamma_{m}} \tag{8.146}
\end{equation*}
$$

where $\left\langle V^{2}\right\rangle$ is the mean square of the random interaction. In the other notations, the nonuniform transformation of the energy scale (8.146) is given by the equation

$$
\begin{equation*}
E=\tilde{E}+\left\langle V^{2}\right\rangle \operatorname{Tr} G(\tilde{E}) \tag{8.147}
\end{equation*}
$$

where $\hat{G}_{0}(\tilde{E})$ is the unperturbed Green's function.
Equation (8.147) implicitly determines the dependence of new energies on the old ones, that is, the function $\tilde{E}(E)$, and thus enables us to calculate the new shapes of any spectral dependence $\tilde{\Phi}(E)$, provided we knew the old profiles $\Phi(E)$. This relation reads

$$
\begin{equation*}
\tilde{\Phi}(E)=\Phi[\tilde{E}(E)] \tag{8.148}
\end{equation*}
$$

Note that the transformation of the spectra, given by (8.147) and (8.148), is not as trivial as it may look at first glance: Plural roots of (8.147) result in the multivalued dependence of $\tilde{E}(E)$. In this case we may obtain a superposition of profiles corresponding to different roots. For the roots of (8.147) to be taken into account, we use the following selection rule. We find the contour in the complex plane $\tilde{E}$ going from $-\infty$ to $\infty$, on which the imaginary part of the right-hand side of (8.147) equals zero. We take only the roots of (8.147) located on this contour, and numerate them according to the sequence in which they are ordered, starting from $-\infty$. The even roots give positive and the odd, negative contributions. Note that the roots of the transcendental equation (8.147) obtained in this way are not necessarily real: they usually have negative imaginary parts.

For the derivation of (8.147), there is no need to perform all of the calculations based on the ensemble averaging of the all-order perturbation series of the Green's functions in the general case, since all of the physically important
ingredients of such a calculation have already been illustrated in Sect.4.1.3 for the particular example of the Green's function projected to the unperturbed energy eigenstate $X_{m}(E)=\langle m|\langle G(E)\rangle|m\rangle$. Equation (4.21) can be directly generalized to the average Green's functions in any representation, which reads

$$
\begin{equation*}
\langle G(E)\rangle=\sum_{m} \frac{|m\rangle\langle m|}{E-\Delta_{m}-\left\langle V^{2}\right\rangle \operatorname{Tr}\langle G(E)\rangle} . \tag{8.149}
\end{equation*}
$$

We simply notice that substitution of the renormalized energy $\widetilde{E}=E-$ $\left\langle V^{2}\right\rangle \operatorname{Tr}\langle G(E)\rangle$ into the definition of the Green's function

$$
\begin{equation*}
G(E)=\sum_{m} \frac{|m\rangle\langle m|}{E-\Delta_{m}} . \tag{8.150}
\end{equation*}
$$

given in the representation of the unperturbed energy eigen states $|m\rangle$, immediately gives $\langle G(E)\rangle=G(\widetilde{E})$, while from the invariance of the trace operation one obtains $\operatorname{Tr}\langle G(E)\rangle=\operatorname{Tr} G(\widetilde{E})$, which results in (8.147).

## Example of Two Merging Multiplets

Let us illustrate this rule by simple examples. For a degenerate multiplet containing $N$ components we find

$$
\begin{equation*}
\tilde{E}-E+\frac{N\left\langle V^{2}\right\rangle}{\tilde{E}}=0 . \tag{8.151}
\end{equation*}
$$

This equation has two roots

$$
\begin{equation*}
\tilde{E}_{1,2}=\frac{E}{2} \pm \sqrt{\left(\frac{E}{2}\right)^{2}-N\left\langle V^{2}\right\rangle} \tag{8.152}
\end{equation*}
$$

The sign selection rule suggests minus, and the new Green's function therefore reads

$$
\begin{equation*}
\langle G(E)\rangle=\frac{N}{\tilde{E}}=\frac{N}{\frac{E}{2}-\sqrt{\left(\frac{E}{2}\right)^{2}-N\left\langle V^{2}\right\rangle}} \tag{8.153}
\end{equation*}
$$

For $\left(\frac{E}{2}\right)^{2}<N\left\langle V^{2}\right\rangle$ the imaginary part of the Green's function yields the Wigner semicircular state density

$$
\begin{equation*}
g(E)=\operatorname{Im} G(E)=\left(\pi\left\langle V^{2}\right\rangle\right)^{-1} \sqrt{E^{2}-4 N\left\langle V^{2}\right\rangle} \tag{8.154}
\end{equation*}
$$

shown in Fig. 4.19.
One can also see the coalescence of two Wigner semicircles. We take $N_{a}$ levels, at the energy $\Delta_{a}=\Delta$ and $N_{b}$ levels of the energy $\Delta_{b}=-\Delta$. The trace of the imaginary part of the Green's function $\operatorname{Im} G(E)=\operatorname{Im}\left[N_{a}(E-\right.$


Fig. 8.10. Density of states $g(E)$ calculated with the help of (8.156) for vareous parameters $V^{2} N / \Delta^{2}$ (a) and $C=N_{a} /\left(N_{a}+N_{b}\right)$ (b). This is an analog of the semicircular Wigner distribution (Fig. 4.19) for the case of two multiply degenerate levels coupled by a random matrix.
$\left.\Delta-i 0)^{-1}+N_{b}(E+\Delta-i 0)^{-1}\right]$ gives us the unperturbed density of states. For the transformed energy $\tilde{E}$ one finds a cubic equation

$$
\begin{equation*}
\tilde{E}-E+V^{2} G(\tilde{E}) \equiv \tilde{E}-E+\frac{V^{2} N_{a}}{\tilde{E}-\Delta-i 0}+\frac{V^{2} N_{b}}{\tilde{E}+\Delta-i 0}=0 \tag{8.155}
\end{equation*}
$$

with the roots $\tilde{E}_{1,2,3}(E)$ given by the Cardano formula. The selection rule yields the root with the negative imaginary part. We note that the variable $E$ is real, and therefore according to (8.149), the density of states $g(E)$ is given by $g(E) \equiv \operatorname{Im} \operatorname{Tr}\langle\hat{G}(E)\rangle=\operatorname{Im} G[\tilde{E}(E)]$. After substitution of $G[\tilde{E}(E)]$ from (8.155) this reads

$$
\begin{equation*}
g(E)=-\operatorname{Im} \tilde{E}_{3}(E) \tag{8.156}
\end{equation*}
$$

In Fig. 8.10 we depict the density of states for various concentrations $N_{a} /\left(N_{a}+N_{b}\right)$.

## Examples of Multidimensional Motion

For a free particle moving in $d$-dimensional space, with kinetic energy proportional to the momentum squared $p^{2}$, the density of states as a function of energy $E$ reads $g(E)=d(2 \pi)^{-d / 2} E^{d / 2-1} / \Gamma(d / 2)$, and therefore the trace of the Green's function in the eigenstate representation reads

$$
\begin{equation*}
\operatorname{Tr} G(E)=\frac{d}{(2 \pi)^{d / 2} \Gamma(d / 2)} \int_{0}^{\infty} \frac{x^{d / 2-1} e^{-\nu x}}{E-x-i 0} d x \tag{8.157}
\end{equation*}
$$

where a small cut-off constant $\nu$ is introduced in order to avoid the divergency of the upper limit. For $d=1,2,3(8.157)$ yields

$$
\begin{align*}
\operatorname{Tr} G(E) & =\frac{i}{\sqrt{2(E-i 0)}} \\
\operatorname{Tr} G(E) & =\frac{1}{\pi}[\ln (-E+i 0)+\ln \nu+\mathbf{C}] \\
\operatorname{Tr} G(E) & =\frac{3 i}{\sqrt{2} \pi} \sqrt{E-i 0}-3 \sqrt{\frac{\pi}{2 \pi^{3} \nu}} \tag{8.158}
\end{align*}
$$

respectively, where $\mathbf{C} \simeq 0.577 \ldots$ is the Euler constant. We do not include here the dimensional parameters related to the mass of the particle and the size of the system, but simply define the mean squared perturbation by the rate $W=\left\langle V^{2}\right\rangle g$ of relaxation that moves the particle with a given momentum out of the momentum state. This implies that the perturbation amplitude is strong and rapidly changing in space, such that this rate is independent of the momentum.

We will not take into account the energy-independent terms of (8.158) that just shift the energy scale without affecting the profile of the density of states, but concentrate only on the energy-dependent ones, that have to be substituted into (8.147). This yields

$$
\begin{align*}
& E=\tilde{E}-\frac{i}{\sqrt{2(\tilde{E}-i 0)}} W \\
& E=\tilde{E}+\frac{W}{\pi}[\ln (-\tilde{E}+i 0)] \\
& E=\tilde{E}-\frac{3 i W}{\sqrt{2} \pi} \sqrt{\tilde{E}-i 0} \tag{8.159}
\end{align*}
$$

where the first and the last equations can be solved analytically, whereas the second one requires a numerical solution. Equations (8.159) have a certain symmetry with respect to scaling and shift operations, such that their solutions result in the universal functional dependencies of the density of states, whereas the energy axis experiences shift and scaling, which may also be accompanied by a scaling of the density axis. Indeed, representing $E$ and $\tilde{E}$ in the first equation as the quantities $X_{1}$ and $\widetilde{X}_{1}$, respectively, that are both multiplied by the factor $W^{2 / 3}$, one gets rid of the parameter $W$. $A$ similar representation of $E$ and $\tilde{E}$ in the second equation containing the factor $W / \pi$, followed by the shift of the argument $X_{2} \rightarrow X_{2}-\ln W / \pi$ gives a similar result. The right-hand side of the third equation is quadratic with respect to $\sqrt{\tilde{E}}$. Therefore, by simple shift of the argument $E$ by the value $-9 W^{2} / 8 \pi^{2}$ followed by taking the roots and squares one finds that the combination $-\operatorname{Im}(\tilde{E}) / \pi W$ remains the same square root function, simply shifted along the energy axis. We therefore arrive at

$$
\begin{align*}
& \widetilde{g}_{1}(E)=\frac{-1}{\pi W^{1 / 3}} \operatorname{Im} \widetilde{X}_{1}\left(\frac{E}{W^{2 / 3}}\right) \\
& \widetilde{g}_{2}(E)=-\frac{1}{\pi} \operatorname{Im} \widetilde{X}_{2}\left(\frac{E \pi}{W}-\ln \frac{W}{\pi}\right) \\
& \widetilde{g}_{3}(E)=\frac{3}{\sqrt{2} \pi^{2}} \sqrt{E-\frac{9 W^{2}}{8 \pi^{2}}} \tag{8.160}
\end{align*}
$$

where $\widetilde{X}_{1}(X)$ and $\widetilde{X}_{1}(X)$ are universal functions given by the solutions of the algebraic and transcendental equations

$$
\begin{gather*}
\left(\widetilde{X}_{1}-X\right)^{2} 2 \widetilde{X}_{1}+1=0  \tag{8.161}\\
X=\widetilde{X}_{2}+\left[\ln \left(-\widetilde{X}_{2}+i 0\right)\right]
\end{gather*}
$$

respectively. The first equation implies a standard solution for the cubic equation, while we have already encountered the second equation in (5.160) in the context of a band interacting with the moving level. There we were interested in the real solution of the latter equation that determines the position of an adiabatic state near the continuum edge, whereas now we are concentrating on the imaginary solution describing the transformation of the state density. In Fig. 8.11 we show the profiles of the imaginary parts of these functions.


Fig. 8.11. Universal profiles of the state density $g(E)$ perturbed by a random matrix as a function of the scaled energy $E$, given by the solutions $\tilde{X}_{1}(X)$ (solid line) and $\pi^{-1} \tilde{X}_{2}(\pi X)$ (dash-dot line) of (8.161) for one- and two-dimensional motion, respectively.

## Examples of Absorption Lines

The density of states is not the only characteristic of a quantum system that is linear in the Green's function. The other important physical quantity is the
transition probability. The natural example of this is the photon absorption cross-section, corresponding to dipole transitions from the ground state of a quantum system to a set of highly excited levels, coupled among themselves by a random perturbation. When the perturbation mixes up the dipole accessible odd states with even states, forbidden for the dipole transitions, the profile $F(E)$ of the absorption line experiences a significant transformation. We illustrate such a transformation for two merging multiplets that undergo transformation given by (8.155). We assume that only one of these multiplets, say $a$, is accessible from the ground state $|0\rangle$ and that the other multiplet is "dark". Therefore the profile

$$
\begin{equation*}
F(E)=\langle 0| \mu_{0 a} G(E) \mu_{a 0}|0\rangle \tag{8.162}
\end{equation*}
$$

given for $\left\langle V^{2}\right\rangle=0$ in terms of the unperturbed Green's function (8.150) and the matrix elements $\mu_{0 a}$ of the dipole moments (taken here as identical) adopts the form

$$
\begin{equation*}
\langle F(E)\rangle=\frac{N_{a}\left|\mu_{a 0}\right|^{2}}{\tilde{E}(E)-\Delta-i 0}, \tag{8.163}
\end{equation*}
$$

as suggested by the transformation (8.149) of the energy scale in the presence of the perturbation.

In Fig. 8.12 we show the evolution of $\langle F(E)\rangle$ for an increasing, random perturbation $\left\langle V^{2}\right\rangle$. For small perturbations the dipole activity is localized


Fig. 8.12. Absorption contour of a "bright" multiplet merging with a "dark" one as a result of random coupling of the multiplet levels, calculated with the help of (8.163). (a) The level scheme for transitions to coupled bright and dark multiplets. (b) The absorption profiles $\sigma(\Delta, V)$ for small, medium and strong random coupling. (c) The same as a function of the mean squared coupling. (d) The same near the point where the multiplets merge.
on the manifold $a$ and gives the "triangular" shaped profile typical of the multiplets with a Wigner semicircular shape. When the nearest edges of the manifolds $a$ and $b$ get closer, part of the dipole activity moves to the multiplet $b$, while the superposition with a non-uniform state density at the edges creates a two-hump profile, shown in Fig. 8.12(b). For strong coupling, when the multiplets merge, the line acquires an almost Lorenzian shape which returns to the "triangular" one only for very strong couplings, when both multiplets can be treated as one.

## Universal Features in Time Dependencies

When considering in Sect. 4.3 the population dynamics of two degenerate levels, we encountered a non exponential behavior of the population difference in the long time limit. The non-exponential time dependence also manifests itself in the frequency domain as discontinuities of derivatives of the corresponding spectra. The sharp edges of the Wigner semicircular distribution entering the imaginary part of $Q \sim \widetilde{E}$ in (4.54) as the origin of the oscillating dependencies given by Bessel functions in (4.61) for the population difference is one of the examples. We can illustrate this by taking the Fourier transform of the Green's function in (8.153) and calculating the response function

$$
\begin{equation*}
R(t) \equiv \int_{\infty}^{-\infty} \frac{d E}{2 \pi} \tilde{G}(E) e^{i \varepsilon t}=\int_{\infty}^{-\infty} \frac{d E}{2 \pi} \frac{e^{i E t}}{E-\sqrt{E^{2}-4 N\left\langle V^{2}\right\rangle}}=\frac{J_{1}\left(t \sqrt{N\left\langle V^{2}\right\rangle}\right)}{t \sqrt{N\left\langle V^{2}\right\rangle}} \tag{8.164}
\end{equation*}
$$

The square of the response function $R(t)$ gives us the time-dependent probability of the transition. We recall the asymptotic behavior of the Bessel function $J_{1}(x) \sim x^{-1 / 2} \sin (x-\pi / 4)$ and see that this probability decreases as $t^{-3}$. In this respect, the Wigner semicircular spectrum behaves like a sixdimensional random walk

We can also find some typical features of the optical response function $R(t)$ corresponding to the profile

$$
\begin{equation*}
\langle F(E)\rangle=\langle 0| \mu_{0 a}\langle G(E)\rangle \mu_{a 0}|0\rangle \tag{8.165}
\end{equation*}
$$

in the general case. Calculating the integral

$$
\begin{equation*}
R(t)=\int_{\infty}^{-\infty} \frac{d E}{2 \pi}\langle 0| \mu_{0 a} G[\tilde{E}(E)] \mu_{a 0}|0\rangle e^{i E t} \tag{8.166}
\end{equation*}
$$

of the inverse Fourier transformation by parts one finds

$$
\begin{equation*}
R(t)=-\int_{\infty}^{-\infty} \frac{e^{i E t}}{2 \pi i t} \frac{\partial\langle 0| \mu_{0 a} G[\tilde{E}(E)] \mu_{a 0}|0\rangle}{\partial E} d E \tag{8.167}
\end{equation*}
$$

A spectrum edge located at a point $E=E_{e d}$ where $G[\tilde{E}(E)] \sim\left(E-E_{e d}\right)^{1 / 2}$ gives the contribution

$$
\begin{equation*}
R(t) \sim \frac{\exp \left[i E_{e d} t\right]}{t^{3 / 2}} \tag{8.168}
\end{equation*}
$$

to the integral (8.167), and hence the overall response in the long-time limit decreases as $t^{-3 / 2}$ and results from the interference of responses (8.168) of all spectrum edges. This yields a multifrequency interference pattern in the case of a spectrum consisting of many spectral bands. However when for a certain value of the random coupling two closest edges of neighboring bands merge at a point $E=E_{\text {cusp }}$, as shown in Fig. 8.10, the density of state form a cusp with the profile $\sim\left|E-E_{\text {cusp }}\right|^{5 / 3}$ and the corresponding contribution to the response

$$
\begin{equation*}
R(t) \sim \frac{\exp \left[i E_{\text {cusp }} t\right]}{t^{8 / 3}} \tag{8.169}
\end{equation*}
$$

decays faster as compared to the edge contributions.

### 8.3.2 The Effect of Quantum Recurrences on the State Density Profiles

Thus far we have been considering the situation when the typical strength $V$ of the random perturbation is much larger than the typical spacing $1 / g$ between the neighboring levels of the unperturbed system and we were concentrating on the transformation of the spectral density on a scale much larger than $1 / g$. In other words, the time uncertainty corresponding to our spectral resolution was much shorter than the typical Heisenberg return time $\sim g$. Now we turn to the case when the perturbation is of the order of the mean spacing $V g \sim 1$ and the time uncertainty of the resolution has to be larger than the return time. The influence of quantum returns therefore becomes dominant.

In Fig. 8.13(a) we illustrate this situation by showing the average density of states $g(\varepsilon, \Delta)$ in the most simple case of a "complex" quantum system consisting of two levels with an energy spacing $2 \Delta$, which is perturbed by a Hermitian $2 \times 2$ random matrix $V_{k m}$. Averageing of the state density

$$
\begin{array}{r}
\quad \delta\left(\varepsilon-\frac{V_{11}+V_{22}}{2}+\sqrt{\left(\frac{V_{11}-V_{22}}{2}+\Delta\right)^{2}+V_{12} V_{21}}\right) \\
+\delta\left(\varepsilon-\frac{V_{11}+V_{22}}{2}-\sqrt{\left(\frac{V_{11}-V_{22}}{2}+\Delta\right)^{2}+V_{12} V_{21}}\right)
\end{array}
$$

with the Wigner perturbation distribution function $g \sim e^{-V_{11}^{2}-V_{22}^{2}-2 V_{12} V_{21}}$ reads

$$
\begin{equation*}
\frac{(\varepsilon+\Delta) e^{2 \Delta \varepsilon}-(\varepsilon-\Delta)}{2 \sqrt{2 \pi} \Delta} \exp \left[-\frac{(\varepsilon+\Delta)^{2}}{2}\right] \tag{8.170}
\end{equation*}
$$

as one can find by performing the direct integration over $d V_{11} d V_{22} d\left|V_{12}\right|$ $d \arg V_{12}$. Here $\varepsilon$ stands for energy which is scaled in the mean perturbation


Fig. 8.13. The density of quantum states at a resolution scale less than the mean level spacing is affected by the interference of recurrences. For the simplest, twolevel "complex" quantum system (a) it results in a two hump profile (8.170). (b) Consistent description of larger systems requires that one takes into account that different trajectories (dash lines) give identical contributions when they correspond to the the same oriented graph (solid lines). (c) Example of an orgraph with parallel transitions. For this orgraph $K_{0}=4 ; K_{1}=3 ; K_{2}=2 ; p_{0,1}=p_{1,0}=2 ; p_{0,2}=$ $p_{2,0}=p_{1,2}=p_{2,1}=1$. The number of different Euler paths is 15 according to (8.175) and, indeed, the orgraph corresponds to the 15 trajectories: 010121030; 010102120; 010120210; 010201210; 010210120; 010212010; 012010210; 012021010; 012101020; 012102010; 020101210; 020121010; 021010120; 021012010; 021201010.
$\left\langle V^{2}\right\rangle^{1 / 2}$ as well as all other energy units in this subsection. One sees the two Gaussian profiles corresponding to two levels that for a vanishing spacing $\Delta$ form a single two hump profile. The humps originate from the level repulsion by the interaction $V_{12}$.

## Trajectories and Oriented Graphs

In order to obtain the state density in the general case, we begin the exact analytical calculation of the averaged perturbation series (4.2) which is quite unwieldy but contains no approximation apart from the ensemble average with the distribution function $g \sim \exp \left(-\operatorname{Tr} \widehat{V}^{2}\right)$. In particular, the calculation does not imply that the diagrams have only a tree-like structure, which was the key assumption in the regime of short times $t<g$. The final result of these calculations is a threefold integral representation of the average Green's function $\langle G(\varepsilon)\rangle$ which gives (8.170) as a particular case.

A generic term of the perturbation series corresponds to a diagram which may return many times to the same level. It has the structure

$$
\begin{equation*}
S\left(\left\{p_{n, m}\right\}\right)=\prod_{n, m}\left[\left(V_{n m}\right)^{p_{n m}}\right] \prod_{n}\left[\frac{1}{\left(\varepsilon-\Delta_{n}\right)^{K_{n}}}\right] \tag{8.171}
\end{equation*}
$$

where $K_{n}=\sum_{m} p_{m n}$ is the number of times the level $|n\rangle$ has been visited, and the numbers $p_{m n}=0,1,2 \ldots$ show how many times the matrix element $V_{n m}$ is found on the trajectory. Each trajectory $0 \rightarrow \ldots \rightarrow n \rightarrow m \rightarrow \ldots \rightarrow 0$ starts and ends at the state $|0\rangle$, the only initially populated level. For the state $|0\rangle$ the power index is $K_{0}=\sum_{m} p_{m 0}+1$. When needed, one obtains the trace of the Green's function, by taking for the state $|0\rangle$, all levels of the system and summing up the resulting contributions. The Green's function is the sum of the amplitudes $S$ of all trajectories

$$
\begin{equation*}
G(\varepsilon)=\sum_{\text {traject. }} S\left(\left\{p_{n m}\right\}\right) \tag{8.172}
\end{equation*}
$$

where many of these trajectories may have the same set $\left\{p_{n, m}\right\}$ of indices and therefore give identical contributions.

Let us represent the states $n$ of the system by points in a plane, and depict a self-intersecting trajectory by connecting these points by the lines representing the coupling $V_{n m}$ which has been encountered on the trajectory. We obtain an oriented graph (orgraph), that is, the totality $\left\{p_{n, m}\right\}$ of lines connecting each pair $n, m$ of the points, and $p_{n m}$ gives the number of different lines going from the point $n$ to the point $m$. The set of numbers $\left\{p_{m n}\right\}$ uniquely defines the orgraph.

There is no one-to-one correspondence between trajectories and orgraphs: Any trajectory of the sum (8.172) corresponds only to one orgraph, but this orgraph may also correspond to some other trajectories. One can say that each of these trajectories goes round the orgraph in its own way, as shown in Fig. 8.13(b). All of the trajectories corresponding to the same orgraph have identical amplitudes $S$ and hence they interfere constructively. Therefore the sum (8.172) of the terms (8.171) can be represented in the form

$$
\begin{equation*}
G(\varepsilon)=\sum_{\operatorname{orgraph}\left\{p_{n, m}\right\}} \prod_{n}\left[\frac{\prod_{m}\left(V_{n m}\right)^{p_{n, m}}}{\left(\varepsilon-\Delta_{n}\right)^{K_{n}}}\right] \mathcal{E}\left(\left\{p_{m n}\right\}\right), \tag{8.173}
\end{equation*}
$$

where the factor $\mathcal{E}\left(\left\{p_{m n}\right\}\right)$ allows for a number of different trajectories on the orgraph $\left\{p_{n m}\right\}$.

## The Number of Euler Paths and Sum Over Carriers

We now make use of the combinatorial expression for $\mathcal{E}\left(\left\{p_{m, n}\right\}\right)$ entering (8.173) which is known explicitly for $p_{n, m}=0,1$. It is the so-called number of Euler paths, given as

$$
\begin{equation*}
\mathcal{E}=\prod_{n}\left[\left(K_{n}-1\right)!\right] \text { minor }\left|\delta_{n, m}\left(\sum_{m} p_{n m}\right)-p_{n m}\right|, \tag{8.174}
\end{equation*}
$$

where the symbol "minor" stands for any of the $M^{2}$ identical $(M-1)$-th order minors corresponding to the $M$ th-order matrix $\delta_{n m} \sum_{m} p_{n m}-p_{n m}$. The determinant $\left|\delta_{n m} \sum_{m} p_{n m}-p_{n m}\right|$ of this matrix equals zero. Instead of calculating minors, one can calculate the determinant where unity is added to the diagonal element corresponding to the level $|0\rangle$. In this way, one can replace the sum of $p_{m n}$ by the value of $K_{n}$. We illustrate this in the orgraphs by adding an entering arrow to the point corresponding to $|0\rangle$. One obtains the generalization of (8.174) for the case $p_{n, m}=0,1,2, \ldots$ dividing it by the number $\prod_{n, m} p_{n, m}$ ! of the transposition of all identical lines in the orgraph

$$
\begin{equation*}
\mathcal{E}\left(\left\{p_{n, m}\right\}\right)=\prod_{n}\left[\frac{\left(K_{n}-1\right)!}{\prod_{m} p_{n m}!}\right]\left|\delta_{n, m} K_{n}-p_{n, m}\right| \tag{8.175}
\end{equation*}
$$

The summation over orgraphs (8.173) is not equivalent to the summation over all $p_{n, m}=0,1, \ldots$, because a trajectory may not visit all of the states of the system. As a consequence, the order of the matrix $\delta_{n, m} K_{n}-p_{n m}$ in (8.175) can be different for different orgraphs. We therefore perform the summation in two steps. Firstly we select the carrier $C(\{n\})$ of the orgraph, that is, a group $\{n\}$ of $N$ levels, each of which must be visited by the corresponding trajectory at least once. The carrier always includes the state $|0\rangle$. Then we sum the contributions of all orgraphs with identical carriers, that is, we perform the summation over all $p_{n, m}$ running from one to infinity. In this sum we keep the size $N \times N$ of the matrix $p_{m, n}$ unchanged. When some of the numbers $K_{n}=$ $\sum_{m} p_{n, m}$ equal zero, the determinant in (8.175) has a zero line and hence the corresponding $\mathcal{E}$ vanishes. Hence the imposed condition, that each trajectory must visit all the states of the carrier, holds automatically. Secondly we sum up the contributions of all possible selections of the carriers with different $N$ and arrive at

$$
\begin{equation*}
G(\varepsilon)=\sum_{C(\{n\})} \sum_{p_{n, m}=0}^{\infty} \prod_{n}\left[\frac{\prod_{m}\left(V_{n m}\right)^{p_{n, m}}}{\left(\varepsilon-\Delta_{n}\right)^{K_{n}}}\right] \mathcal{E}\left(\left\{p_{m n}\right\}\right) . \tag{8.176}
\end{equation*}
$$

Note that if we want to keep the natural enumeration of the states included in a carrier $(n, m=1,2, \ldots)$, we have to attribute to them new numbers according to the sequence of the initial numeration. The other possibility is to retain the original enumeration, which does not run over all natural numbers. It is convenient to choose the set-theory notation that gives identical expressions for both types of enumeration. We write $l, m \in\{n\}$ if the states $l$ and $m$ are included in the set $\{n\}$, and $l, m \in\{n\} /\{q\}$ if at the same time these states are not included in the set $\{q\}$. All states in the sets are ordered according to the original enumeration.

## Factorization of $\mathcal{E}$ and the Ensemble Average

Our aim now is to factorize the number of Euler paths and represent the sum of products $\sum_{p_{n m}} \prod_{n m}$ as the product of sums $\prod_{n m} \sum_{p_{n m}}$. Following that, we perform the summation for each pair $n, m$ independently. We employ the relation

$$
\begin{equation*}
\frac{\left(K_{n}-1\right)!}{\left(\varepsilon-\Delta_{n}\right)^{K_{n}}}=\int_{C_{t}} \frac{d t_{n}}{t_{n}} t_{n}^{K_{n}} \exp \left\{-\left(\varepsilon-\Delta_{n}\right) t_{n}\right\} \tag{8.177}
\end{equation*}
$$

were $C_{t}$ is the integration contour that starts at $t=0$ and goes to infinity along the direction ensuring convergence of the integrals, and thereby in the integrand we obtain the product of the factors $t_{n}^{p_{m n}}$. Separation of the determinants can be performed with the help of Grassmann integration. Here we employ another approach, which to a certain extent illustrates the meaning of the graded anticommuting variables. We note that $\frac{\partial}{\partial \alpha} \exp \left(p_{l m} \alpha\right) \rightarrow p_{l m}$ for $\alpha \rightarrow 0$ and write

$$
\begin{align*}
& \left|\delta_{l m} \sum_{k \in\{n\}} p_{k m}+\delta_{l 0} \delta_{m 0}-p_{l m}\right| \\
= & \lim _{\alpha \rightarrow 0}\left|\delta_{l m} \sum_{k \in\{n\}} \frac{\partial}{\partial \alpha_{k m}}+\delta_{l 0} \delta_{m 0}-\frac{\partial}{\partial \alpha_{l m}}\right| \exp \sum_{l, m} p_{l m} \alpha_{l m}, \tag{8.178}
\end{align*}
$$

where $\delta_{l 0} \delta_{m 0}$ allows for the fact that $K_{0}=\sum_{m} p_{0 m}+1$.
Let us make parallels between the technique based on the differential operator-determinant and the technique based on graded variables. Interchange of two neighboring columns in a determinant results in a change of its sign. This property corresponds to the anticommutation relations of the Grassmann variables $\chi \chi^{\prime}=-\chi^{\prime} \chi$, whereas the fact that elements of each line or column enter into each term of the determinant only once along with the requirement $\alpha \rightarrow 0$ corresponds to the vanishing of all powers of these variables. The fact that the argument should depend on the variables $\alpha_{l m}$ for Grassmann variables $\chi$ corresponds to the condition $\int d \chi=0$.

We substitute Eqs. (8.177), (8.178), and (8.175) into (8.176) and arrive at

$$
\begin{align*}
G(\varepsilon)= & \lim _{\alpha \rightarrow 0} \sum_{C(\{n\})}\left|\delta_{l m} \sum_{k \in\{n\}} \frac{\partial}{\partial \alpha_{k m}}+\delta_{l 0} \delta_{m 0}-\frac{\partial}{\partial \alpha_{l m}}\right|  \tag{8.179}\\
& \prod_{l \in\{n\}} \int_{C_{t}} \frac{d t_{t}}{t_{l}} t_{l}^{\delta_{l l}} \exp \left\{-\left(\varepsilon-\Delta_{l}\right) t_{l} \prod_{m \in\{n\}}\left[\sum_{p_{n, m}=0}^{\infty} \frac{\left[t_{l} \exp \left(\alpha_{l m}\right) V_{l, m}\right]^{p_{l m}}}{p_{l, m}!}\right] .\right.
\end{align*}
$$

Here by the power $\delta_{l 0}$ we take into account that all carriers $C(\{n\})$ include the state $n=0$. In the summation over $p_{l, m}=0,1, \ldots$ we recognize the Taylor expansion of the exponents, and hence

$$
\begin{align*}
G(\varepsilon)= & \lim _{\alpha \rightarrow 0} \sum_{C(\{n\})}\left|\delta_{l m} \sum_{k \in\{n\}} \frac{\partial}{\partial \alpha_{k m}}+\delta_{l 0} \delta_{m 0}-\frac{\partial}{\partial \alpha_{l m}}\right|  \tag{8.180}\\
& \prod_{l \in\{n\}} \int_{C_{t}} \frac{d t_{l}}{t_{l}} t_{l}^{\delta_{l}} \exp \left\{-\left(\varepsilon-\Delta_{l}\right) t_{l} \prod_{m \in\{n\}}\left[\exp \left(t_{l} V_{l, m} \exp \alpha_{l m}\right)\right] .\right.
\end{align*}
$$

Now we perform the average over the Gaussian distribution of the matrix elements $V_{n m}$. Application of (8.102) to each matrix element $V_{l m}$ and its complex conjugate counterpart $V_{m l}$ yields

$$
\begin{equation*}
\left\langle\prod_{l, m \in\{n\}} \exp \left(t_{l} V_{l, m} \exp \alpha_{l m}\right)\right\rangle=\exp \left[\sum_{l, m \in\{n\}} t_{l} t_{m} \exp \left(\alpha_{m l}+\alpha_{l m}\right)\right] \tag{8.181}
\end{equation*}
$$

We now move the differential operator-determinant under the integral (8.180) and substitute (8.181) such that the integrand of the many-fold integral over $\prod d t_{l}$ adopts the form

$$
\begin{align*}
\mathcal{Y}= & \lim _{\alpha \rightarrow 0}\left|\delta_{l m} \sum_{k \in\{n\}} \frac{\partial}{\partial \alpha_{k m}}+\delta_{l 0} \delta_{m 0}-\frac{\partial}{\partial \alpha_{l m}}\right|  \tag{8.182}\\
& \exp \left[\sum_{l, m \in\{n\}} t_{l} t_{m} e^{\alpha_{m l}+\alpha_{l m}}\right] \prod_{l \in\{n\}}\left\{t_{l}^{\delta_{l, 1}-1} \exp \left[-\left(\varepsilon-\Delta_{l}\right) t_{l}\right]\right\} .
\end{align*}
$$

Now we calculate the derivatives, and take the limit $\alpha=0$. The determinant in (8.182) is a differential operator. It comprises only the first-order derivatives over each particular $\alpha_{l, m}$, since all of the terms generated by a determinant include only the matrix elements corresponding to different lines. Therefore in the last sum we can expand the exponents

$$
t_{l} t_{m} e^{\alpha_{m l}+\alpha_{l m}}=t_{l} t_{m}+t_{l} t_{m} \alpha_{m l}+t_{l} t_{m} \alpha_{m l} \alpha_{l m}
$$

move the zero-order terms in front of the determinant, and keep only the terms linear in $\alpha$. Bilinear terms are not important in the limit $\alpha \rightarrow 0$. Indeed, when we take a bilinear term $\alpha_{m l} \alpha_{l m}$ and consider the structure of the minor of the determinant comprising derivatives over both these variables, we find that it has the structure

$$
\left|\begin{array}{cc}
\frac{\partial}{\partial \alpha_{m l}} & -\frac{\partial}{\partial \alpha_{m l}} \\
-\frac{\partial}{\partial \alpha_{l m}} & \frac{\partial}{\partial \alpha_{l m}}
\end{array}\right|=0
$$

vanishing identically. The integrand of the averaged equation (8.180) given by (8.182) adopts the form

$$
\begin{align*}
\mathcal{Y}= & \left|\delta_{l m} t_{m} \sum_{k \in\{n\}} t_{k}+\delta_{l 0} \delta_{m 0}-t_{l} t_{m}\right| \\
& \exp \left[\sum_{l, m \in\{n\}} t_{l} t_{m}\right] \prod_{l \in\{n\}}\left\{t_{l}^{\delta_{l, 1}-1} \exp \left[-\left(\varepsilon-\Delta_{l}\right) t_{l}\right]\right\} \tag{8.183}
\end{align*}
$$

The determinant can be found explicitly, and for the $N \times N$ matrix reads $\left(\sum_{l \in\{n\}} t_{l}\right)^{N-2} \prod_{l \in\{n\}} t_{l}$. We can represent this in completely factorized form with the help of the integrals

$$
\begin{align*}
\left(\sum_{l \in\{n\}} t_{l}\right)^{-1} & =\int_{0}^{\infty} \exp \left\{-u \sum_{l \in\{n\}} t_{l}\right\} d u  \tag{8.184}\\
\left(\sum_{l \in\{n\}} t_{l}\right)^{N-1} & =\frac{1}{2 \pi i} \int_{C_{0}} \frac{d x}{x} \int_{0}^{\infty} d z e^{-z}\left(\frac{z}{x}\right)^{N-1} \exp \left\{-x \sum_{l \in\{n\}} t_{l}\right\}
\end{align*}
$$

where the contour $C_{0}$ circumvents the point $x=0$ and allows one to select the ( $N-1$ )-th power of the sum according the Cauchy formula for derivatives of analytical functions, and the integral over $d z$ compensates for the factorial in the Cauchy formula. For the factorization of the sums in the exponent, we employ a simpler version

$$
\begin{equation*}
\exp \left(\sum_{l \in\{n\}} t_{l} \sum_{m \in\{n\}} t_{m}\right)=\frac{1}{\sqrt{\pi}} \int_{-\infty}^{\infty} \exp \left(-y^{2}-2 y \sum_{m \in\{n\}} t_{m}\right) d y \tag{8.185}
\end{equation*}
$$

of the standard factorization procedure (8.102). Substituting (8.184) and (8.185) into (8.183) for the integrand we obtain (8.180) in the form

$$
\begin{align*}
\langle G(\varepsilon)\rangle= & \sum_{C(\{n\})} \frac{1}{2 \pi i \sqrt{\pi}} \int_{C_{0}} \frac{d x}{x} \int_{0}^{\infty} d z \int_{-\infty}^{\infty} d y \int_{0}^{\infty} d u\left(\frac{z}{x}\right)^{N-1} e^{-z-y^{2}} \\
& \prod_{l \in\{n\}} \int_{C_{t}} t_{l}^{\delta_{0 l}} \exp \left[-\left(u+x+2 y+\varepsilon-\Delta_{l}\right) t_{l}\right] d t_{l} \tag{8.186}
\end{align*}
$$

The last integration is standard and yields $\left(u+x+2 y+\varepsilon-\Delta_{l}\right)^{-1-\delta_{0 l}}$ since the integration contour $C_{t}$ can be directed from $t_{l}=0$ to $-i \infty$ employing an infinitesimal imaginary part of the variable $\varepsilon$.

Now we perform the summation over all of the carriers that contain the state $n=0$. All other states of the system can either be included or not included in the carrier $C(\{n\})$. We attribute the factor unity to the states which are not included in the carriers, and extend our product to all states of the system. Then the summation over all possible carriers reduces to the multiplication of a binomial, and we arrive at

$$
\begin{align*}
\langle G(\varepsilon)\rangle= & \frac{1}{2 \pi i \sqrt{\pi}} \int_{C_{0}} \frac{d x}{x} \int_{0}^{\infty} d z \int_{-\infty}^{\infty} d y e^{-z-y^{2}} \int_{0}^{\infty} d u  \tag{8.187}\\
& \left(u+x+2 y+\varepsilon-\Delta_{0}\right)^{-2} \prod_{l=1}^{N}\left[1+\frac{z}{x}\left(u+x+2 y+\varepsilon-\Delta_{l}\right)^{-1}\right] .
\end{align*}
$$

Equation (8.187) corresponds by construction to the matrix element $G(\varepsilon)_{0,0}$ of the Green's operator, since we have considered only the trajectories that start at the level $|0\rangle$. In order to obtain the trace of the Green's operator we note that the factor in (8.187) corresponding to the starting level $|0\rangle$ is the derivative $\partial / \partial u$ of the factor identical to that of the other levels, and therefore after integration over $d u$ the trace reads

$$
\begin{align*}
\operatorname{Tr}\langle G(\varepsilon)\rangle= & \frac{1}{2 \pi i \sqrt{\pi}} \int_{C_{0}} \frac{d x}{x} \int_{0}^{\infty} d z \int_{-\infty}^{\infty} d y e^{-z-y^{2}} \\
& \frac{x}{z}\left\{\prod_{l=0}^{N}\left[1+\frac{z}{x}\left(x+2 y+\varepsilon-\Delta_{l}\right)^{-1}\right]-1\right\} \tag{8.188}
\end{align*}
$$

where the integration constant -1 is chosen in order to exclude the nonphysical constant term of the product, corresponding to the carrier containing no level.

Further simplifications have no physical background, being purely technical. We employ an integral representation similar to (8.185), replacing $e^{-y^{2}}$ by the integral of $\pi^{-1 / 2} \int e^{-q^{2}-2 i y q} d q$, change the variables $2 y+x+\varepsilon \rightarrow 2 y$ and $x \rightarrow z x$, simultaneously also changing the order of the integration when needed, perform integration over $d z$, and arrive at

$$
\begin{align*}
\operatorname{Tr}\langle G(\varepsilon)\rangle= & \frac{1}{2 \pi^{2} i} \int_{-\infty}^{\infty} d q \int_{C_{0}} \frac{d x}{x} \int_{-\infty}^{\infty} d y e^{-q^{2}-2 i y q+i q(x z+\varepsilon)} \\
& \frac{1}{1-i q x} x\left\{\prod_{l=0}^{N}\left[1+\frac{1}{x}\left(2 y-\Delta_{l}\right)^{-1}\right]-1\right\} . \tag{8.189}
\end{align*}
$$

The integration contour $C_{0}$ is a loop separating the high-order poles at the point $x=0$ from the first-order pole at $x=1 / i q$. We move this contour to infinity, where the integral vanishes and calculate the residual at the point $x=1 / i q$. This yields

$$
\begin{align*}
\operatorname{Tr}\langle G(\varepsilon)\rangle= & \frac{1}{\pi} \int_{-\infty}^{\infty} d q \int_{-\infty}^{\infty} d y e^{-q^{2}-2 i y q+i q \varepsilon} \\
& \frac{1}{i q}\left\{\prod_{l=0}^{N}\left[1+i q\left(2 y-\Delta_{l}\right)^{-1}\right]-1\right\} \tag{8.190}
\end{align*}
$$

which is the final result for the general case. The limit of a large number of states follows from (8.190) after some algebraic manipulations, when we assume that each factor of the product is close to unity and set

$$
\prod_{l=0}^{N}\left[1+i q\left(2 y-\Delta_{l}\right)^{-1}\right]=\exp \left[i q \sum_{l=0}^{N}\left(2 y-\Delta_{l}\right)^{-1}\right]
$$

by analogy to (3.118) and (8.111).


Fig. 8.14. The density of quantum states for multiplets perturbed by a random matrix. The even number of states $N=2,4,6,8$ are shown by solid lines, while the dashed lines show the profiles for the odd $N$. The asymptotic semicircular dependence, scaled to the parameters of $N=10$, is also shown. The two-hump dependence of $N=2$ coincides with that of Fig. 8.12.

The other interesting limit is that of identical levels, that is $\Delta_{l}=0$ for all $l$. In this case one can expand the binomial $\left(1+i q(2 y)^{-1}\right)^{N}$ and integrate each term independently allowing for the integral representation

$$
H_{n}(x)=\frac{2^{n}}{\sqrt{\pi}} \int_{-\infty}^{\infty}(x+i \widetilde{q})^{n} e^{-\widetilde{q}^{2}} d \widetilde{q}
$$

for the Hermite polynomials, which yields

$$
\begin{equation*}
\langle G(\varepsilon)\rangle=\frac{e^{-\varepsilon^{2} / 4}}{2 \sqrt{\pi}} \sum_{k=0}^{N-1} \frac{N!H_{2 k}(\varepsilon / 2)}{2^{k}(k+1)!k!(N-k-1)!} \tag{8.191}
\end{equation*}
$$

In Fig. 8.14 we show the density of states for the case of $N=1,2, \ldots 8$ along with the Wigner semicircular profile corresponding to the asymptotic of $N \rightarrow$ $\infty$.

We conclude by summarizing the main results of this subsection. We sum up explicitly the all-order perturbation series for an arbitrary quantum system perturbed by a random matrix from a Gaussian unitary ensemble. We take into account all self-intersections of the trajectories exactly, with the help of the explicit expression for the number of Euler paths. We obtain the exact expression for the transformed spectra in the form of a two-fold integral. A simple transformation rule, which governs the "macroscopic" transformation of spectra, follows from the exact expression in the extreme of a large number of states.

### 8.3.3 The Density of Quantum States of Fractals

The description of the interference phenomena inhibiting the level's decay and affecting the state density at high resolution was based on the analysis of self-intersections of the transition trajectories going over quantum states of a complex quantum system. Now we focus on another case of interference, when the orgraphs of the diagrams do not self-intersect but correspond to a tree-like topological structure. However each branch of the "tree" may be passed an infinite number of times by the trajectory. One encounters a physical realization of this situation by considering a quantum particle moving along a ramified structure like that shown in Fig. 2.20. In contrast with the problems considered earlier in the energy representation, for this problem the coordinate representation is much more natural, whereas the Green's function of the particle is given by the Feynman path integral in its classical formulation, that is as a sum over all possible trajectories, weighted by a phase factor suggested by the classical action along each trajectory.

## The Renormalization Equation

In Fig. 8.15 we compare and contrast the tree-like diagrams for both cases, assuming that each vertex is connected to only two vertexes placed at a larger distance from the beginning of the tree. The Green's function of the particle

$$
\begin{equation*}
G_{n m}(p)=\exp \left(-i p L_{m n}\right) \tag{8.192}
\end{equation*}
$$

moving with momentum $p$ along a one-dimensional interval of length $L_{m n}$ connecting nodes $n$ and $m$ of the ramified structure, corresponds to the line of the graph which coincides with standard notations in solid state physics and field theory. The vertexes that were denoting the Green's functions $\left(\varepsilon-\Delta_{k}\right)^{-1}$ of the unperturbed complex system correspond now, on the contrary, to the interaction matrix

$$
\widehat{S}=\left(\begin{array}{ccc}
S_{b b} & S_{b r} & S_{b l}  \tag{8.193}\\
S_{r b} & S_{r r} & S_{r l} \\
S_{l b} & S_{r r} & S_{r r}
\end{array}\right)
$$

which describes the scattering of the particle at a node of the ramified structure. We assume that the scattering amplitudes are identical for all vertexes, depending only on the direction. The subscripts $r$ and $l$ denote the right and left direction, respectively, while $b$ stands for the direction towards the root of the tree. We also assume that the distance $L_{m n}$ can decrease by a factor $a$ with each additional step in the number of intervals $N$ separating the interval from the root. For simplicity, we take all intervals corresponding to the same $N$ to be identical, that is $L_{m n}=L_{N}$.

Let us derive an equation for the Green's function (8.192) which would be an analog of the renormalization equations (4.21) or (8.149) for complex systems. The derivation is illustrated in Fig. 8.16(a). By $X_{N}$ we denote the


Fig. 8.15. Tree-like diagrams for the probability amplitude at times shorter than the recurrency time (left) and the diagrams for a particle moving along a ramified structure (right) have the same topological structure. However line and vortex symbols have opposite meanings. The free propagation of the complex system corresponds to the Green's function $\left(\varepsilon-\Delta_{k}\right)^{-1}$ shown by the vertex point, which corresponds to the line $G_{n, m}=\exp \left(-i p L_{n, m}\right)$ for the propagating particle. The coupling $V_{n, m}$ corresponds in turn to the vertex $\hat{S}$ describing scattering of the particle in the node of the fractal. The length of the intervals $L_{n, m}$ may decrease with increasing number $N$ of steps separating a vertex from the root.
renormalized Green's function for a particle leaving a node $N$ for the neighboring vertexes $N+1$ and coming back. This implies the relation $\phi_{r}=X_{N} \psi_{r}$ between the amplitude of the returning wave $\phi_{r}$ and the amplitude $\psi_{r}$ of the wave leaving a node $N$ along the right channel. The same condition for the incoming $\phi_{l}$ and outgoing $\psi_{l}$ waves holds for the left channel. The renormalized Green's function $X_{N-1}$ is a product of three factors, two functions $G_{N-1}$ acquired before and after the passinge the vertex $N$ and the ratio $\psi_{b} / \phi_{b}$ allowing for the scattering process in the node. By the relation $\psi=\widehat{S} \phi$, the scattering matrix (8.193) relates the vector $\phi=\left(\phi_{b}, \phi_{r}, \phi_{l}\right)$ of the incoming amplitudes for the vortex $N$ with the vector $\boldsymbol{\psi}=\left(\psi_{b}, \psi_{r}, \psi_{l}\right)$ of the outgoing amplitudes.

For further analysis, we take the simplest model of the unitary scattering matrix $\widehat{S}$, completely symmetric with respect to the interchanging channels, which also satisfies the long-wave requirement, giving the $-\pi$ phase shift for the amplitudes $\phi_{b} / \psi_{b}=-1$ when the same is valid for two other pairs $\phi_{r} / \psi_{r}=-1$ and $\phi_{l} / \psi_{l}=-1$. The latter condition means that if both channels $r$ and $l$ are closed, channel $b$ sees the vertex as an infinite, vertical, completely reflecting wall. In other words, no quantum defect is associated


Fig. 8.16. Recurrence relation for the Green's functions. (a) The relations among the scattering amplitudes at a node $N$ have to be complemented by the propagation conditions given by the phase factors $X_{N}$. (b) The trajectory $X\left(p / a^{N}\right)$ of mapping (8.196) shows the value of the asymptotic dependence after each consecutive momentum scaling. Starting from the vicinity of the unstable stationary point with $X=-1$ which corresponds to the boundary condition $\psi\left(x_{N}\right)=0$ for the wavefunction in the coordinate representation, the trajectory goes to the stable stationary point that corresponds to the condition $\psi_{x}^{\prime}\left(x_{N}\right)=0$ of the constant flux (c).
with the vertex. The only matrix satisfying these properties gives the following system of equations

$$
\widehat{S} \boldsymbol{\phi}=\left(\begin{array}{ccc}
-1 / 3 & 2 / 3 & 2 / 3 \\
2 / 3 & -1 / 3 & 2 / 3 \\
2 / 3 & 2 / 3 & -1 / 3
\end{array}\right)\left(\begin{array}{c}
\phi_{b} \\
X_{N} \\
\psi_{r} \\
X_{N} \\
\psi_{l}
\end{array}\right)=\left(\begin{array}{c}
\psi_{b} \\
\psi_{r} \\
\psi_{l}
\end{array}\right)
$$

and after excluding $\psi_{r}$ and $\psi_{l}$ via the relation $X_{N-1}=G_{N-1}^{2} \psi_{b} / \phi_{b}$ it yields the renormalization condition

$$
\begin{equation*}
X_{N-1}\left(L_{N-1} p\right)=G_{N-1}^{2} \frac{1-3 X_{N}\left(L_{N} p\right)}{X_{N}\left(L_{N} p\right)-3} \tag{8.194}
\end{equation*}
$$

relating two successive renormalized Green's functions $X$.

## Asymptotic Ddependence

With the help of (8.194), one recursively finds the renormalized Greens function on a larger scale, provided they are known for the closest smaller scale, and vice versa. In the long wave limit, when $G_{N-1}^{2}=\exp \left(-i p L_{N-1}\right) \simeq 1$ this equation is a mapping which does not contain any parameters depending on the scale size, apart from the arguments of $X$. Each branch of the tree-like diagram has the same topology. If we assume that the size of the branches scales
by a factor $1 / a$ when we pass from the branch $N-1$ to the next branch $N$, the corresponding functions $X_{N-1}\left(L_{N-1} p\right)$ and $X_{N}\left(L_{N} p\right)$ should be identical, differing only by the scale factor of the argument, that is $X_{N-1}\left(L_{N-1} p\right)=$ $X\left(L_{N-1} p\right)$ and $X_{N}\left(L_{N} p\right)=X\left(L_{N} p\right)=X_{N-1}\left(L_{N-1} p / a\right)$. In other words, for small $p$ (8.194) reads

$$
\begin{equation*}
X(p)=\frac{1-3 X(p / a)}{X(p / a)-3} ; \quad \text { or } \quad X(p / a)=\frac{1+3 X(p)}{X(p)+3} \tag{8.195}
\end{equation*}
$$

This implies that starting from the value of the asymptotic dependence $X(p)$ for a given argument $p$ and performing the mapping

$$
\begin{equation*}
X_{N}=\frac{1+3 X_{N-1}}{X_{N-1}+3}, \quad p=p / a \tag{8.196}
\end{equation*}
$$

one can obtain the values of this function at all points $p=p / a^{N}$ approaching the point $p=0$, as shown in Fig. 8.16(b).

The mapping (8.196) has two stationary points at $X= \pm 1$, as one can see in Fig. $8.16(\mathrm{~b})$. The unstable stationary point $X=-1$ corresponds to the reflection of the wavefunction from the node as from a rigid potential wall, where the probability amplitude $\psi_{b}+\phi_{b}$ must tend to 0 . Therefore the reflected wave $\psi_{b}$ is shifted in phase by $\pi$ with respect to the incident wave $\phi_{b}$. The stable stationary point $X=1$ corresponds to the probability amplitude flux going through the node to the ensuing branches of the fractal structure and returning. The total length of the branches can be very large, even infinite, and hence the totality of the branches occupies a vast domain of the phase space and is therefore capable of absorbing a big fraction of the population. However, the length of each particular trajectory going through the node to the end of the ramified structure decreases, and often even vanishes at $N \rightarrow \infty$, such that the total accumulated phase of the returning wave does not have enough distance to change. This implies that in this case, the wavefunction at the node satisfies the boundary condition with zero derivative $p \psi_{b}-p \phi_{b}=0$, and the reflection occurs with zero phase shift, as illustrated in Fig. 8.16(c). The latter condition changes the spectrum of the system drastically, since now the Bohr quantification rule permits momenta values close to zero. Let us call "the spectrum of fractals" the part of the spectrum at small $p$ resulting from the large total length of the ramified structure and consider it in more detail.

The asymptotic behavior of $X_{a s}(p)$ for $p \rightarrow 0$ can be found exactly. Comparing the mapping (8.196) with the trigonometric identity

$$
\tanh (z+\ln \sqrt{2})=\frac{1+3 \tanh z}{3 \tanh z+3}
$$

we find that $X_{a s}(p)=\tanh [-\alpha \ln (L p)]$ satisfies (8.195) exactly, provided that $\alpha \ln (a)=\ln \sqrt{2}$. Therefore

$$
\begin{equation*}
X_{a s}(p)=\frac{1-(L p)^{2 \alpha}}{1+(L p)^{2 \alpha}} \tag{8.197}
\end{equation*}
$$

which indeed tends to -1 for $p \rightarrow 0$. Starting with $X_{a s}(p)$ for a small argument $p / a^{N_{0}}$ and performing $N_{0}$ iterations back to the larger scales, one finds the exact dependence $X(p)$ for all $p$ with any predetermined accuracy, depending on the number $N_{0}$ of iterations. Even for a few iterations, the result perfectly reproduces all of the important features of the system, whereas $N_{0}=10$ is sufficient for practical reasons for $a>1.5$ and gives an accuracy better than $1 \%$. Note that the asymptotic form (8.197) remains valid even in the case when the scaling factor $a$ depends on the interaction number $N$, with the only difference that we have to substitute the average of $\ln (a)$ over all of the iterations performed into the definition $\alpha \ln (a)=\ln \sqrt{2}$ of the index $\alpha$. Note that $\alpha$ coincides with the fractal dimensionality of the tree-like structure.

## Density of States

In order to find the state density profile for the fractal structures, we note that at a given energy two different momenta with the same absolute values $|p|$ and opposite signs contribute to the state density. The corresponding Green's functions interfere in the coordinate representation. Therefore, instead of generalization of the expression $g(E)=-\operatorname{Im}[\operatorname{Tr} G(E) / \pi] \mathrm{n}$ the case of the double degeneracy of the energy levels, it is more convenient to make use of the known Fabri-Perrot formula

$$
\begin{equation*}
g(|p|)=\frac{1}{2 \pi} \operatorname{Re} \frac{1+X(p)}{1-X(p)} \tag{8.198}
\end{equation*}
$$

which for $X(p)=R \exp [i L p]$ gives the state density of the rectangular potential well, with the reflection coefficient $R$ at the edges. Actually, apart from a constant factor, this formula coincides with the expression for the transmission profile of the Fabri-Perrot resonator for the particle wave. In Fig. 8.17(a) we show the state density as a function of $|p|$, calculated for different scaling factors $a$ with the help of (8.198). For $a>1$ but relatively small, the position of the first resonance goes down with respect to the position of the first state in the potential well, which is consistent with the increase of the space domain occupied by the particle. The presence of the fractal manifests itself in the broadening of the profile. For larger $a$ the resonance returns to the FabriPerrot position, and goes even higher, being repelled by the states appearing at the domain $|p| \rightarrow 0$. The appearance of these state close to zero energy is the most important manifestation of the fractal structure, and implies that the population is getting distributed over all of the domain occupied by the fractal.

The asymptotic dependence $g(|p|) \sim(L p)^{-2 \alpha}$ directly follows from the substitution of (8.197) into (8.198). In Fig. 8.17(b) we show this domain of


Fig. 8.17. Spectra of fractal structures. (a) Density $g$ of states as a function of dimensionless momentum $L p$. Fabri-Perot resonances (dotted line) corresponding to the isolated first branch get shifted towards lower $L p$ when additional branches with a relatively small scaling factor ( $a=1.7$ solid line) are added. For a larger scaling factor ( $a=2.9$ dot-dash line and $a=3.9$ dashed line) the resonances are broadened and shifted due to the presence of states at small $L p$. (b) The states at small $L p$ typical of fractals for the same parameters. The state density $g(L p)$ follows power laws with different indices for different scaling factors. (c) Fractal state density transformed by a random perturbation for different fractal dimensionality found from the solution of (8.199) with $\beta=2 \alpha$. Dashed lines show the particular cases corresponding to the Wigner semicircle and to one- and two-dimensional free particles, depicted in Figs. 4.19, 8.14, and 8.11, respectively. (The state density $g$ is positive in both upward and downward directions.)
momenta with a higher resolution, where the power dependencies different for different $a$ are clearly seen. We note that the relation between the energy and momentum can be different for different cases, changing from $E \sim p^{2}$ for free particles to $E \sim p$ for the conducting electrons near the Fermi surface in metals. In the latter case the energy dependence of the state density for small $E$ reads $g(E) \sim E^{-\ln 2 / \ln a}$. For this case, an important qualitative conclusion can be made, - for $a>2$ the total number of low-laying states is finite, since the integral $\int_{0} g(E) d E$ converges at the lower limit. This is consistent with the fact that the total length of the fractal branches is finite in this case, and is given by the sum of the geometric series $L /(1-2 / a)$. For $a<2$ the total length of the branches is infinite, and the integral diverges. For a free particle with $E \sim p^{2}$ the density of states at small energies behaves as
$g(E) \sim E^{-\alpha-1 / 2}$, and therefore the combination $2(1-\alpha)$ can be interpreted as a spatial dimensionality of the fractal, by comparison with the expression $g(E) \sim E^{(d-2) / 2}$ for the regular $d$-dimensional space.

## Randomly Perturbed Fractals

We now address the question: how general are our results obtained for the specific scattering at the nodes and for the strict scaling law of the branch sizes? One possibility to answer the question is to consider deviations of branch lengths from the mean value $L / a^{N}$ and scattering matrices from the mean matrix $S$ of (8.193) as random perturbations. This assumption allows us to employ the transformation $(8.147)$, $(8.148)$ that can be written in the dimensionless variables

$$
\begin{equation*}
E=\tilde{E}-(-\tilde{E}+i 0)^{-2 \alpha} \tag{8.199}
\end{equation*}
$$

by analogy to (8.159) for a free particle moving in regular space, whereas the mean squared random perturbation just scales the energies. This yields universal shapes of the state densities for each value of the fractal dimension shown in Fig. 8.17(c). Note that this form is identical for the reciprocal dimensions.

## 9 Bibliography and Problems

Even considering particular aspects of problems that have been encountered through many decades in different branches of physics and chemistry, it is impossible to give a comprehensive review of all relevant publications. This is particularly the case for the dynamics of multilevel quantum systems, which usually is just a technical aspect and not the main concern of research papers and textbooks dedicated mostly to physical phenomena, rather than to the technical details of their description. In considering a bibliographic search in such a situation, one has to make a difficult choice amongst the numerous possibilities. The main objective of the bibliography presented here was to find a balance between the pioneering papers, topical reviews, monographs, textbooks concerning particular physical domains, and the original publications most closely related to the dynamical processes in complex quantum systems. This review therefore consists of two parts: the "Bookshelf" comprising a selection of classical original papers and monographs, and the "Problems" where a selection of more recent publications from different domains is suggested to the reader for consideration from the point of view of the dynamics of multilevel systems developed in the previous chapters. Though some old classical papers cited below might be more difficult to access than their subsequent interpretations, we believe that the pleasure of following the development of the original ideas of the great physicists should definitely compensate for this trouble.

Among numerous modern publications, only a limited number of papers from a number of different fields have been selected. The majority of the problems considered there have been formulated in terms traditional to the corresponding physical domains and we suggest that they be approached from the point of view of the dynamics of complex quantum systems, since we believe that the effort aimed at such a diversification of approaches will help the reader to attain a deeper insight into the nature of these processes. Consideration of some of the problems requires a certain advanced experience in ajacent fields, and that can be achieved by consulting the cited monographs and reviews. In particular, this relates to the field-theoretical models in combination with the saddle-point approximation that have not been considered here in sufficient detail.

### 9.1 Bookshelf

## Cascades and Exponential Decay

"Consider an ensemble of two-level systems interacting with a quantized oscillator of a field mode in a highly exited Fock state in the presence of an infinite number of other field modes in the vacuum state" - that is how we would formulate today the problem considered by A. Einstein[73] in his famous paper Zur Quantentheorie der Strahlung (1917). Solving this problem, we would draw the level scheme, Fig. 8.1(d), consider the corresponding Schrödinger equation, and find that the system is described by the master equation with the probability amplitudes $B n$ coupling the closest states of the excited field oscillator and the atomic ensemble, and the probabilities $A$ coupling states of the atomic ensemble with the bands of states describing the reservoir of the other field oscillators. Of course Einstein did not formulate this problem in terms of the Schrödinger equation or Dirac formalism that were introduced later, but this problem can still be viewed as the first example of the dynamics of complex quantum systems. This approach was first adopted a decade later by L. D. Landau[123] in his paper Das Dämpfungsproblem in der Wellenmechanik (1927)

The description of the interaction between radiation and matter in terms of multilevel quantum systems is already found in the paper Berechnung der natürlichen Linienbreite auf Grund der Diracschen Lichttheorie (1930) by V. Weisskopf and E. Wigner[175], dedicated to the analogy between the width of the resonance of a classical harmonic oscillator with the friction and the spectral line widths. The dynamical aproach had been employed by E. Fermi[80] in his paper Quantum theory of radiation (1932), where the level-band problem had been formulated and solved in the context of atomic spontaneous emission. The radiative widths of resonances were the main physical quantities of interest. Later on he generalized the result and obtained the "Fermi golden rule" for the transition rate in an arbitrary level-band system, valid for times shorter than the Heisenberg time.

The dynamical consideration of multilevel quantum systems for a long time remained an auxiliary tool in the investigation of the spectral properties of quantum objects that manifest themselves in scattering processes. G. Breit and E. Wigner[28] in their paper Capture of slow neutrons (1936) had demonstrated that nuclear interactions can result in scattering profiles much broader than those suggested by the radiation life-times. The development of the microwave technique had offered an opportunity to control the strength of interaction between radiation and matter, such that the rates of radiative transitions became of the order of the interactions in quantum systems. R. Karplus and J. Schwinger[110] in their paper A note on saturation in microwave spectroscopy (1948) have shown that the width of the resonance can depend on the radiation field intensity.

## Alternative Line Shapes

For open quantum systems, there always exists a time parameter which gives the typical time of an exponential decay. However some important spectral features may result from another, non-exponential type of dynamics, which yields different line shapes. One of the most striking theoretical results is the semicircular profile of the energy eigenstate density obtained by E . Wigner[177] in his paper On the distribution of the roots of certain symmetric matrices (1957). This distribution emerges from the model of a Gaussian random matrix based on the central limit theorem of statistics. Though the physical validity of this model for the description of complex nuclei has been questioned by F. Dyson[69] in his series of papers Statistical theory of the energy levels of complex systems (1962), it nevertheless remains an elegant example of a complex system being adequate and useful in other domains of physics, in situations where the spectral complexity truly results from the randomness of the system parameters.

Corrections to the semicircular Wigner distributions near the spectrum edges had been explored by Ya. G. Sinai and A. B. Soshnikov[171] in their paper A refinement of Wigner's semi-circle law in a neighborhood of the spectrum edge for random symmetric matrices (1998), while the general properties of the Hamiltonians perturbed by a random matrix had been discussed by E.Brezin, S. Hikami, and A.Zee[33] in the paper Universal correlations for deterministic plus random Hamiltonians (1995).

The other famous example of a spectral profile different from a Lorentzian one is given by U. Fano[76] in the paper Effects of configuration interaction on intensities and phase shifts (1961) which was a developpment of the ideas formulated in his earlier publication[75] Sullo spettro di assorbimento dei gas nobili presso il limite dello spettro d'arco (1935). Formulated initially for purely intra-atomic interactions, this example became of extreme importance with the development of lasers in the situations when laser-induced interactions strongly perturb atomic spectra. The dynamics corresponding to these spectra also differ from a simple exponential decay, as illustrated by P. L. Knight, M. A. Lauder and B. J. Dalton[111] in the review paper Laser-induced continuum structure (1990).

Non-Lorentzian spectral profiles are typical of solid state physics. J. Hubbard[96] in his series of works entiteled Electron correlations in narrow energy bands (1963-1964), had demonstrated that under certain conditions, spectra of random alloys have a semicircular structure resulting from electron correlation. The Hubbard model Hamiltonian for the electron liquid does not have any random parameters, however the resulting spectra so resemble those obtained via the random perturbation of a system of two degenerate levels, shown in Fig. 8.10, that it invites one to think about the fundamental reasons for such a similiarity.

## Dynamics

The general relation between the dynamical and spectral properties of quantum systems had been established by R. Kubo[118] in his paper A general expression for the conductivity tensor (1956), which gives the linear response of the system to a weak perturbation. This expression implies that the system is subjected to a relaxation, such that the nonlinear effects do not have time to manifest themselves.

Hamiltonian systems behave differently. The first exemple of such behavior had been given by E. L. Hahn[95] in his paper Spin Echoes (1950), where it had been shown that an ensemble of spins in a slightly inhomogeneous field experience a coherent damping which can be reversed by sending a $\pi$-pulse which inverts the phases of the probability amplitudes. A very convinient tool for the consideration of this kind of problem had been introduced in the paper by I. I. Rabi, N. F. Ramsey, and J. Schwinger[154] entitled Use of rotating coordinates in magnetic resonance problems (1954).

The interation of the ensemble of two-level systems with radiation in the regime of developed coherent dynamics also differs considerably from the linear responce regime. R. H. Dicke[60] in his paper Coherence in spontaneous radiation processes (1954) had demonstrated that the spontaneous decay rate resulting from the interaction of the ensemble with a photon vacuum increases strongly for the highly correlated states of individual atoms achieved in the course of their coherent dynamics. These states also change the propagation of the electromagnetic radiation in a medium composed of the ensemble of two-level atoms as shown by C. K. Rhodes, A. Szöke, and A. Javan[156] in the paper The influence of level degeneracy on the self-induced transparency effect (1968), and by S. L. McCall and E. L. Hahn[135] in the paper Selfinduced transparency (1969).

## Localization and Chains

The dynamics of complex quantum systems are strongly related to the structure of the energy eigenfunctions. P. W. Anderson[7] in his famous paper Absence of diffusion in certain random lattices (1957) had demonstrated that in a perturbed infinite complex system with random parameters, the diffusion process can be partially or completly suppressed, such that the energy eigenfunctions turn out to be localized on a limited number of unperturbed states. In the one-dimensional case, the corresponding energy eigenfunctions are localized exponentially, as had been numerically demonstrated in the paper by R. E. Borland[27] entitled The nature of the electronic states in disordered one-dimensional systems (1963). An analytically soluble model had been introduced by P. Lloyd[131], in his paper Exactly solvable model of electronic states in a three-dimensional disordered Hamiltonian: non-existence of localized states (1969). D. J. Thouless[174] in his paper Maximum metallic resistance in thin wires (1977) had discussed the localization conditions for the one-dimensional case.

Exhaustive numerical analysis of the one-dimensional case had been presented in the paper by K. Ishii[100] entitled Localization of eigenstates and transport phenomena in the one-dimensional disordered systems (1972) whereas the analytical consideration had been summarized in the monograph by I.M. Lifshits, S.A. Gredeskul and L.A. Pastur[129] Introduction to the theory of disordered systems (1988). It turns out that the localization phenomenon is not something specific to disordered systems, but has deeper roots in the algebraic structure of operators, as had been shown in the paper by Ya. G. Sinai[169] entitled Anderson localization for one-dimensional difference Schroedinger operator with quasi-periodic potential (1987).

Quantum dynamics in chains had been considered for the first time in the paper by V. L. Berezinskii[15] entitled Kinetics of a quantum particle in a one-dimensional random potential (1974). This work had demonstrated that the all-order correlations of the Green function are important in the onedimensional case, contrary to the Bogolyubov method for the description of complex dynamics.

## Ensembles. Levels Statistics

The idea of a statistical description of complex quantum systems had first been employed by E. Wigner[176] in his paper On a class of analytic functions from the quantum theory of collisions (1951) for the neutron scattering amplitudes that have been assumed to possess certain statistical properties in the complex plane of energies. Much later, F. Dyson[69] in his paper Statistical theory of the energy levels of complex systems, (1962) had proposed the ensemble average as a universal approach to complex quantum systems, which was relevent even in the case where no real physical ensemble had been associated with the complex system, whereas the statistical approach emerged from the "renounce of the exact knowladge... of the nature of the system...". This very concept called, following F. Dyson, "a new kind of statistical mechanics" remains the cornerstone of the modern approach to complex quantum systems.

The dynamical aspects of the behavior of multilevel systems had been restricted to the statistics of resonance widths in open multilevel quantum systems, that had been related in the paper by C.E. Porter and R.G. Thomas [152] entitled Fluctuations of nuclear reaction widths (1956) to the statistics of eigenvectors of complex systems. One finds more details of the statistical approach to nuclear and atomic spectra in the book by C. E. Porter[151] Statistical theories of spectra: fluctuations (1965). The mathemetical background of the statistical approach to complex spectra along with the fundamental results of the level correlation statistics is given in the book by M. L. Metha[136] Random matrices and the statistical theory of energy levels (1967).

In contemporary research, statistical models are widely employed in condensed matter physics. One can consult the review by C. W. J. Beenakker[14] Random matrix theory of quantum transport (1997), the book by J. Imry[99]

Introduction to mesoscopic physics (2002), and the special issue of The Journal of Physics[106] (2003) for a modern view on the statistical description of complex quantum systems in the context of mesoscopic condensed matter physics.

## Classical and Quantum Chaos, Level Statistics

Quantum complexity and chaos have not only a long history of parallel development and mutual enrichment via the exchange of basic ideas but also have common roots in ergodic theory. The first demonstration of the relevance of the statistical approach to the description of dynamical systems had been performed by Ya. G. Sinai[167] in his pioneering paper On the concept of entropy for a dynamic system (1959), that got further development in the later monographs[168] Introduction to ergodic theory (1976), and Topics in ergodic theory (1994)[170] by the same author. About the same time, B. V. Chirikov[46] in his paper Resonance processes in magnetic traps (1960), had introduced the resonance overlap criterion necessary for the destruction of the separability of the motion, which results in the chaotic classical dynamics of mechanical systems. These ideas were further developed in detail in the monographs by A.J. Lichtenberg, M.A. Liberman[128] entitled Regular and stochastic motion (1982), and by G. M. Zaslavski, R. Z. Sagdeev, D. A. Usikov, and A. A. Chernikov[180] entitled Weak chaos and quasi-regular patterns (1991).

The quantum mechanical counterpart of classicaly regular and chaotic systems had been studied numerically. For regular dynamics, the Poissonian statistics of distances to the closest neighboring level and level clustering had been found in the paper by M. V. Bery and M. Tabor[16], entitled Level clustering in the regular spectrum (1977), whereas the level repulsion results were typical of chaotic dynamics, as had been shown in the well-known publication of O. Bohigas, M. J. Giannoni, and C. Schmit[24] entitled Characterization of chaotic quantum spectra and universality of level fluctuation laws (1984). For more details concerning the manifestation of the chaotic dynamics at the quantum level of consideration, one can consult the reviews by B. Eckhardt[72] entitled Quantum mechanics of classically non-integrable systems (1988), and F. M. Izrailev[102], Simple models of quantum chaos: spectrum and eigenfunctions (1990) where some dynamical aspects related to Anderson localization in kicked rotors had been considered, as well as the monographs by M. Gutzwiller[93] Chaos in classical and quantum mechanics (1990), by L. E. Reichl[155] The transition to chaos (1992), and by F. Haake[94] Quantum signatures of chaos (2001) for a general point of view of this field.

An interesting example of level statistics had been discussed in the papers by E. B. Bogomolny, U. Gerland, and C. Schmit[25] entitled Models of intermediate spectral statistics (1999), and[26] Singular statistics (2001), where it had been demonstrated that for essentially quantum systems that have no classical analogs, one could observe intermediate statistics, corresponding
neither to a Poissonian nor to a Wigner-Dyson distribution on the closest neighboring levels spacing.

## Supersymmetry, Unification Disorder and Chaos

Considerable progress in the understanding and a better description of disordered quantum systems has been achieved after the application of field theory methods to disordered quantum systems. This supersymmetry technique, along with the representative references to earlier publications, had been presented in detail in the book by K. Efetov[74] entitled Supersymmetry in disorder and chaos (1997). One should append to this exhaustive publication the review by M. Zirnbauer[183] Anderson localization and non linear sigma-model with graded symmetry (1986) and the recent paper[32] New correlation functions for random matrices and integrals over supergroups (2003) by E. Brezin and S. Hikami.

In the paper Quantum billiards with surface scattering: Ballistic sigmamodel approach (1998) by by Ya. M. Blanter, A. D. Mirlin, and B. A. Muzykantskii[19], the supersymmetry methods have been generalized to the case of ballistically propagating particles which allowd one to apply field theory methods to the ballistic billiards typically considered in the models of quantum chaos.

## Dynamics of Atoms and Molecules

Contrary to the case of nuclear collisions, for the collision of atoms one has the small Born-Oppenheimer parameter $\Lambda=(m / M)^{1 / 4}$ which allows for the fact that the electrons of mass $m$ move much faster relative to the nuclei of mass $M$. The adiabaticity parameter $\Lambda$ is widely employed for the description of atomic collisions, as one can see from the classical monograph by H.S.W. Massey and N. F. Mott[133] entitled The theory of atomic collisions (1997) and the monograph by I. I. Sobelman, L. A. Vainshtain, and E. A. Yukov[172] entitled Excitation of atoms and broadening of spectral lines (1981). When the trajectories of the center-of-mass motion of colliding atoms or molecules can be considered as given, the electron Hamiltonian acquires an explicit time dependence. Numerous models of such collisions have been considered in the monograph[145] Theory elementary atomic and molecular processes (1985) by E. E. Nikitin and in the monograph[146] Theory of slow atomic collisions (1984) by E. E. Nikitin and S. Ya. Umanskii.

Among the pionieering papers related to molecular dynamics, one should cite the paper by A. M. Dykhne[67], entitled Adiabatic perturbation of the descrete spectrum states (1962), where the general semiclassical expression for the transition probability had been obtained, the paper by Yu. N. Demkov and V. I. Osherov[57] entitled Stationary and nonstationary problems in quantum mechanics that can be solved by means of contour integration (1968) where the famous exemple of an exactly soluble problem had been suggested,
and the papers by M. Bixon and J. Jortner[17] entitled Intramolecular radiationless transitions (1968) and[18] Long radiative lifetimes of small molecules (1969) where the multilevel quantum model had been introduced for the description of intramolecular conversion and spontaneous emission.

Complex multilevel systems corresponding to the Rydberg states of atoms and small molecules had been exhaustively described in the monographs by E. W. Schlag[160] entitled ZEKE spectroscopy (1998), the book[107] Molecular applications of quantum defect theory (1996), edited by Ch. Jungen, and the monograph by T. F. Gallagher[87] entitled Rydberg atoms (1994).

## Laser Theory, Quantum Optics

Though Albert Einstein in his work[73] (1917) did not consider the situation where the population of the exited atomic states exceeds the population of the ground state, the master equation derived by him had been employed for the prediction and explanation of induced coherent emission of masers and lasers. Later on it became clear that the master equation did not include information relating to the phase of the electromagnetic field and other, more general, characteristics, such as the photon statistics of the emitted radiation. A consistent description of the interaction of quantized electromagnetic fields with matter had been presented in the review by S. Stenholm[173] entitled Quantum theory of electromagnetic fields interacting with atoms and molecules (1973), where the dynamical description of the quantum evolution of multilevel systems composed of atomic ensembles and multimode fields, arose naturally and described all of the dynamical details of the process, including irreversible relaxation. For more details, one can consult the book by M. Sargent, W. E. Lamb, and M. O. Scully[158] Laser physics (1977) whereas for the modern view on the problems of quantum optics the monograph W. Schleich[161] Quantum optics in phase space (2001) can be recommended.

## Atoms and Molecules in Strong Fields

Atoms and molecules interacting with strong laser fields are the classical example of multilevel quantum systems that experience coherent quantum evolution. F. V. Bunkin and A. M. Prokhorov[35] in their paper The exitation and ionization of atoms in a strong radiation field (1964) had employed the rotating wave aproach, earlier introduced by I. I. Rabi, N. F. Ramsey, and J. Schwinger in the already cited paper[154]. It was the first application of the seminal concept of quasienergy, which was introduced some years later by Ya. B. Zeldovich[181] in his papers The quasienergy of a quantummechanical system subjected to a periodic action (1967) and Scattering and emission of a quantum system in a strong electromagnetic wave (1973), and by V. I. Ritus[157] in the paper Shift and splitting of atomic energy levels by the field of an electromagnetic wave (1967), that had yielded the method of the effective Hamiltonian describing the time evolution of a multilevel quantum system under a strong time-dependent perturbation. The mathematical
formalism of this approach relied on the classical paper by G. Floquet[81] Sur les équations différentialles linéaires à coefficients périodiques (1889). For the modern applications of the Floquet theory to dynamics of atoms in strong fields, one can consult the review by S. Guérin and H. R. Jauslin[92] entitled Control of quantum sdynamics by laser pulses: adiabatic Floquet theory (2003). The dynamics of molecules in laser fields had first been considered by G. A. Askaryan [10] in the paper Excitation and dissociation of molecules in an intense light field (1964) and by F. V. Bunkin, A. M. Prokhorov, and R. V. Karapetyan[36] in the paper Dissociation of molecules in strong radiation field (1965).

Fundamental aspects of the interaction between intense radiation and matter are widely presented in monographs and textbooks. One can consult the book[23] Nonlinear optics by N. Blombergen (1965) for the details of the coupling of photons in atomic media, the book by L. Allen, J. H. Eberly[5] entitled Optical resonance and two-level atoms (1975) for fundamentals of the interaction between strong radiation and matter, the monograph by N.B. Delone, V.P. Krainov[54] entitled Atoms in strong light fields (1985) for a review of multiphoton processes in atoms, and the textbook by C. CohenTannoudji, J. Dupont-Roc, and G. Grynberg[50] Atom-photon interactions: basic processes and applications (1998). The consistent description of atoms and molecules in strong laser fields in terms of multilevel quantum systems is also given in the textbook by N. V. Karlov and V.M. Akulin[4] Intense resonant interactions in quantum electronics (1991), which evidently is a publication that is conceptually closest to the present text.

Interatomic Forces in Resonators and Laser Fields, Ionic Traps, Atomic Optics

The development of lasers in the 1990s had opened up new possibilities for the observation of quantum phenomena. On the one hand, the photon measurement technique had reached the sensibility level of a single photon that had allowed one to address questions about atomic interaction mediated by a single photon in resonators, and the manifestations of such interactions in the dynamics of individual photons. One of the examples of such systems had been presented in the paper by G. Kurizki, A. G. Kofman, and V. Yudson[119] entitled Resonant photon exchange by atom pairs in high- $Q$ cavities (1996). In this context important role played by photonic band structures had been described in detail in the monograph[105] by J. D. Joannopoulos, R. D. Meade, and J. N. Winn Photonic crystals (1995). The novel aspect of the problem was the coupling of internal atomic degrees of freedom with translational atomic motion.

On the other hand, the powerful and stable sources of electromagnetic radiation had enabled experimentalists to control and confine the translational motion of atoms and atomic ions. In the situation where the quantum levels corresponding to the coupled internal and translational motion were
addressed individualy, the multilevel quantum system had been proposed for quantum computations by J. I. Cirac and P. Zoller[48] in their pioneering paper Quantum computation with cold trapped ions (1995). Note that the quantization of the translational motion of atomic ions results in this case from the effective Kapitsa-Dirac potential imposed by rapidly oscillating radiofrequency fields, as had been described in detail in the monograph[89] Ion traps by P. K. Ghosh (1997). This quantization of ionic vibrational modes in ionic traps suggests the comparison of the ensemble of trapped ions with a polyatomic molecule, as has been done in Sect. 2.5.1.

The quantization of the translational motion of neutral atoms can also be achieved by strong laser fields. The fundamentals of such processes are presented in the monograph[137] by P. Meystre entitled Atom optics (2001). The consideration of the aforesaid systems by the methods developed in the previous sections for complex multilevel systems is a matter for future research. For such systems, the complex translational quantum dynamics of atoms results from the spatial structure of the laser field modes.

## Long-Term Dynamics

The main difference between descrete and continuous complex spectra manifests itself in the long-term behavior of the populations. Coherent returns of the population to the initial state were discussed for the first time by J. Eberly, N. B. Narozhny, and J. J. Sanchez-Mondregon [71] in their pioneering paper Periodic spontaneous collapse and revivals in a simple quantum model (1980). Later on, I. Sh. Averbukh and N. F. Perelman[11] in their cornerstone paper Fractional revivals: universality in the long-term evolution of quantum wave packets beyond the correspondence principle dynamics (1989) had elucidated the role of anharmonicity on the violation of the Erenfest principle, and had revealed a rich domain of possible intermediate assymptotics of quantum mechanics resulting from strong and long-lasting interference phenomena in multilevel systems.

The manifestation of such interference in the dynamical properties of multilevel relay-like systems had been numerically shown in the exemple of a hydrogen atom by G. Casati, B. V. Chirikov, and D. L. Shepelyansky[41] in the paper Quantum limitation for chaotic excitation of the hydrogen atom in a monochromatic field (1984) and further discussed by D. L. Shepelyansky[164] in his paper Localization of diffusive excitation in multi-level systems (1987) where the relations between Anderson localization and the long-time assymptotic behavior of complex systems had been demonstrated.

The long-time limit of the coherent dynamics of a strongly perturbed continuum was the main concern of the paper by M. V. Fedorov, M. Yu. Ivanov and A. M. Movsesian[77] entitled Strong-field photoionization of an initially excited hydrogen atom: formation of Rydberg wavepacket, its structure and trapping of population at Rydberg levels (1990), where it had been demonstrated that the life time of the atomic states increases with the increase of
the ionizing laser field intensity. This conterintuitive result is a natural conseqence of correlations between the interaction matrix elements, typical of atomic systems.

In the context of polyatomic molecules, the non-exponential, long-time behavior of complex quantum systems had been discussed by S. A. Schofield and P. G. Wolynes[162] in the paper A scaling perspective on quantum energy flow in molecules (1992), while for metals, a similar approach had been develloped by J. T. Chalker, I. V. Lerner and R. A. Smith[45] in the paper Random walks through the ensemble: linking spectral statistics with wave-function correlations in disordered metals (1996).

## Control

The possibility to control a multilevel quantum system by aplying a timedependent perturbation is based on the existance theorem proved by V. Jurdevic and H. J. Sussmann[108] in their paper Control systems on Lie groups (1972), where it had been shown that the compactness of the group generated by the unperturbed Hamiltonian and the perturbations plays a crucial role and permits one to exert complete control if the group coincides with the group of all possible unitary transformations of the system. Application of this general result to quantum systems has been discussed by A. G. Butkovskiy and Yu. I. Samoilenko[38] in the monograph Control of quantum-mechanical processes and systems (1990). Practical methods of control are presented in the monograph by P. W. Brumer, M. Shapiro[34] entitled Principles of the quantum control of molecular processes (2003). For basic principles of error corrections one can consult the momograph by R. G. Gallager [86] Information theory and reliable communications (1968).

### 9.2 Problems

## Two Levels

1. Describe the dynamics of spin classically (F. Bloch [20], Nuclear induction (1946)) and in terms of two-level systems and Bloch vectors (F. Bloch[21], Generalized theory of relaxation (1956)).

Level-Band

1. Consider the level-band problem in terms of scattering theory. (G. Breit and E. Wigner[28] Capture of slow neutrons (1936)).
2. Consider the level-band model of intramolecular transitions and find the conversion rate of the electronic excitation to vibrational energy (M. Bixon and J. Jortner[17] Intramolecular radiationless transitions (1968)).
3. Consider the transition rates in atoms as a function of the frequency of the probe radiation in the presence of a strong laser field.(P. L. Knight, M. A. Lauder and B. J. Dalton[111] Laser-induced continuum structure (1990)).
4. Consider the interaction of two atoms via a narrow resonator mode (G. Kurizki, A. G. Kofman, and V. Yudson[119] Resonant photon exchange by atom pairs in high- $Q$ cavities (1996)) and generalize the result to the case of a photonic band.
5. Consider the population dynamics corresponding to the Fano spectral profile (Ph. Durand, I. Paladrova and F. X. Cadea[66] Theory of Fano profiles (2001)).
6. Consider the time evolution of an atomic doublet coupled to an ionization continuum (Yu. V. Dubrovskii, M. V. Fedorov, and M. Yu. Ivanov[65] Rydberg atom ionization by an intense laser pulse with smooth time envelope: A model of two closely spaced levels coupled to continuum (1992)). Compare this to the method of projections (L. Mower[141] Decay theory of closely coupled unstable states (1966)).
7. Consider the dynamics of level decay to a band edge, with the desity of states increasing according to the power law. (A. G. Kofman [112] Theory of single-photon bound-free transitions: Extension of the pole approximation (1997)).

## Revivals Wavepackets

1.Condider the regime of recurrences in the molecular quasicontinuum (E. Kyrölä and J. Eberly[121], Quasicontinuum effects in molecular exciation (1985)).
2. Consider the long-time assymptotics of the quasicontinuum and continuous band models (J. Javanainen and E. Kyrölä[104] Long-time limit of a quasicontinuum model (1985)).
3. Consider the time dependence of the fractional revivals spikes (C. Lechtle, I. Sh. Averbukh, and W. P. Schleich[125] Generic structure of multilevel quantum beats (1996)).
4. Consider a band of atomic Rydberg states in terms of wavepackets (M. V. Fedorov, M. Yu. Ivanov and A. M. Movsesian[77] Strong-field photoionization of an initially excited hydrogen atom: Formation of Rydberg wavepacket, its structure and trapping of population at Rydberg levels (1990)).
5. Consider the excitation of an atom by a laser field and its second harmonics in terms of a multilevel system. Interpret the phase dependence of the numerically observed structure ( K. J. Schafer and K. C. Kulander [159] Phase dependent effects in multiphoton ionization induced by a laser field and its second harmonic (1992)).
6. Consider recurrences in a spectrum with the Poisson statistics of levels (P. Pechukas[148] Quantum chaos in the irregular spectrum (1982)).

1. Consider the singular billiard (P. Seba[163] Wave chaos in singular quantum billiard (1990)) and describe it in terms of a level interacting with a band. Compare this to the system perturbed by a rank-1 matrix (E. B. Bogomolny, U. Gerland, and C. Schmit [26] Singular statistics (2001)), and relate the singular statistics to the non-exponential decay of the level to a continuum.
2. Consider a perturbation imposed by a magnetic flux line on a band corresponding to a billiard (R. Narevich, R. E. Prange, O. Zaitsev[144] Square billiard with an Aharonov-Bohm flux line (2001)). Formulate this problem in terms of a rank-1 perturbation.
3. Relate the spectral correlations at the mobility edge of a disordered metal (V. E. Kravtsov, I. V. Lerner, B. L. Altshuler and A. G. Aronov[117] Universal spectral correlations at the mobility edge (1994)) with the powerlaw decay law (J. T. Chalker, I. V. Lerner and R. A. Smith[45] Random walks through the ensemble: Linking spectral statistics with wave-function correlations in disordered metals (1996)).
4. Consider the vibrational dynamics of polyatomic molecules in the regime of power law decay (M. Gruebele[91] Intramolecular vibrational dephasing obays a power law at intermediate times (1998)).

## Two Bands

1. Consider an example of an individual addresing the transitions in a twoband system (I. J. Cirac and P. Zoller[48] Quantum computing with cold trapped ions (1995)).
2. Compare the coherent damping and revivals in the Jannes-Cumming model (J. H. Eberly, N. B. Narozhny, and J. J. Sanchez-Mondregon[71], Periodic spontaneous collapse and revival in a simple quantum model (1980)) with that for an ensemble of spins subjected to a phase conjugating action (E. L. Hahn[95], Spin Echoes (1950)).
3. Consider the two-band model of metallic wires (D. J. Thouless [174], Maximum metallic resistance in thin wires (1977)).
4. Consider the two-band model of unimolecular reactions in laser fields. (M. Quack[153] Theory of unimolecular reactions induced by monochromatic infrared radiation (1978)). Identify the master equation regime, the role of correlations of the dipole matrix elements, and consider the case of several privileged transitions.
5. Determine the probability of multiphoton transition as a function of laser intensity in the presence of removed degeneracy over the momentum projections (V. M. Akulin[2] Contribution of weak and multiphoton transitions to the excitation of polyatomic molecules in an intense laser field (1982)).
6. Consider a Rydberg atom in a strong laser field as a discrete band coupled to a continuum (M. Yu. Ivanov[101] Suppression of resonant multiphoton ionization via Rydberg states (1994)). Demonstrate the role of the
rank-1 interaction matrix in the narrowing of the decay resonances and as an inhibitor of the decay rates. Determine the positions of the quasienergy eigenstates in the complex plane for this problem.

## Quasienergy

1. Derive the time-independent Shcrödinger equation for an atom in a strong, circularly polarized laser field (F. V. Bunkin, A. M. Prokhorov [35] The excitation and ionization of atoms in a strong radiation field (1964)).
2. Consider the role of photon quantum states in the concept of the quasienergy (Ya. B. Zeldovich[181, 182] The quasienergy of a quantummechanical system subjected to a periodic action (1967) and Scattering and emission of a quantum system in a strong electromagnetic wave (1973)).
3. Derive the Schrödinger equation for an atom in an arbitrary periodic external field (M. Gavrila and J. Z. Kaminski[88] Free-free transitions in intense high-frequency laser fields (1984)).
4. Consider the dynamics of Rydberg wave packets in the representation of quasienergy eigenstates (D. Delande and J. Zakrzewski[51] Spontaneous emission of nondispersive Rydberg wave packets (1998)).

## Time-dependent systems

1. Consider a quantum particle scattered by the potrntial $-U_{0} / \cosh ^{2}(x)$ (Landau and Lifshits Quantum mechanics, problem for $\S 23$ ) and trace the analogy with the two-level system excited by the pulse interaction n.5.2.4.
2. Reproduce the result of the paper (Yu. N. Demkov[56] Charge transfer at small resonance defects (1964)) for the exponentially rising coupling between two detuned levels.
3. Consider two coupled levels with a detuning changing parabollically in time (E. E. Nikitin and S. Ya. Umanskii[146] Theory of slow atomic collisions $\S 8.3$ (1984)). Compare the result with the approximate one given by the Dykhne formula (A. M. Dykhne[67] Adiabatic perturbation of the discrete spectrum states (1962)).
4. Estimate the probability of a non-adiabatic transition corresponding to the level-band system of (5.169) with the help of the Dykhne formula.

## Time-Dependent Complex Systems

1. Consider the role of the degeneracy of band states in the DemkovOsherov problem (V. A. Yurovsky, A. Ben-Reuven, P. S. Julienne, and Y. B. Brand[179] Conterintuitive transitions in multistate curve crossing involving linear potentials (1999)).
2. Consider two-state curve crossing of a band (V. A. Yurovsky, A. BenReuven[178] Saturated transition in exactly soluble models of two-state curve crossing with time-dependent potentials (1999)).
3. Estimate the rate of population diffusion over the energy scale for a complex system, adiabatically perturbed by a strong time-dependent perturbation.
4. In the quasienergy representation, consider the decay of a single level to a continuous band of finite width, for the case of harmonic modulation of the coupling (G. S. Agarwal[1] Control of decoherence and relaxation by frequency modulation of a heat bath (2000)).
5. Perform the same for the case of phase modulation (A. G. Kofman and G. Kurizki[114] Universal dynamical control of quantum mechanical decay: Modulation of the coupling to the continuum (2001), A. G. Kofman[113] Relaxation of a two-level system strongly coupled to a reservoir: Anomalously slow decoherence (2001)).
6. Compare the universal population distribution in a quantum system perurbed by a random matrix for colliding molecules (V.M. Akulin[3] Singularities of collisional vibrational relaxation processes in strongly excited spectrally complex polyatomic molecules (1983)) and for quantum dots in an external field (V. E. Kravtsov[116], Time-reversal symmetry breaking by ac field: Effect of commensurability in frequency domain (2002)).
7. Consider a quantum particle evolving in the presence of a chaotic background (D. Kuznetsov, A. Bulgak and G. D. Dang[120] Quantum Lévy process and fractional kinetics (1999)).

## Control

1. Determine a way to perform a given transformation via controlled collision (D. Jaksch, H. Briegel, J. I. Cirac, C. W. Gardiner, and P. Zoller [103] Entanglement of atoms via cold controlled collisions (1999)).
2. Trace the equivalence of the sufficient criterion of the bracket generation condition and the requirement given in the book by P. W. Brumer, M. Shapiro [34] entitled Principles of the quantum control of molecular processes (2003).
3. Check the possibility to control a couple of dipole-dipole interacting spins by application of an external time-dependent magnetic field.

## Group Symmetry Underlying Exact Solubility

1. Show that the system of $n$ two level atoms which has $2^{n}$ quantum levels can be separated into subsystems, each of which is a high-dimensional representation of a $\mathrm{SU}(2)$ group (F. Bloch and I. I. Rabi[22] Atoms in variable magnetic fields (1945)).
2. Consider a harmonic oscilator with time-dependent parameters (K. Husimi [98] Miscellanea in elementary quantum mechanics (1953)). Relate the possibility to find an exact solution with the algebraic properties of the operators entering the Hamiltonian.
3. Consider the underlying group symmetry for the Morse oscillator (R. D. Levine [126] Representation of one-dimensional motion in a Morse potential by a quadratic Hamiltonian (1983)).
4. Generalize the time-dependent oscillator problem to a multidimensional case (V. V. Dodonov and V. I. Man'ko[62] Phase space eigenfunctions of multidimensional quadratic hamiltonians (1986)).
5. Consider the underlying group symmetry for the oscillator with a $1 / x$ perturbation (V. V. Dodonov, V. I. Manko, and D. E. Nikonov[63] Exact propagators for time-dependent Coulomb, delta and other potentials (1992)), and for the Landau levels (V. A. Kostelecky, V. I. Manko, M. M. Nieto, and D. R. Traux[115] Supersymmetry and a time-dependent Landau system (1993)).

## Relay of Isolated Levels

1. Derive the relay-like Schrödinger equation for an ensemble of two-level atoms interacting with the quantized field mode for one-photon and multiphoton cases (S. Stenholm[173] Quantum theory of electromagnetic fields interacting with atoms and molecules (1973)).
2. Perform the same derivation for the photon down-conversion, that is for a field mode nonlinearly coupled to another field mode at half the frequency. (G. Drobny, I. Jex, and V.Buzek [64] Mode entanglement in nondegenerate down-conversion with quantized pump (1993)).
3. Consider an exactly soluble relay and determine the role of relay truncation on the recurrencies dynamics (M. Lindberg and E. Kyrölä[130] Infinite multilevel system coupled to a quantum-field mode (1992)).
4. Find the analytic description of the previous problem. ( B. W. Shore and J. H. Eberly[166] Analytic approximations in multi-level excitation theory (1978)).
5. Determine the population distribution evolution in a perfict lattice with a single defect (P. Lloyd[131] Exactly solvable model of electronic states in a three-dimensional disordered Hamiltonian: non-existence of localized states (1969)).
6. Consider the soluble, semi-infinite relay of levels with a coupling decreasing with the level number and with an arbitrary coupling for the first transition (A. Makarov[132] Coherent excitation of equidistant multilevel systems in a resonant monochromatic field (1977)).
7. Consider the relay with quadratically changing detunings (S. Mukamel and J. Jortner[143] A model for isotope separation via molecular multiphoton photodissociation (1976)) and compare the numerical results of the paper with the predictions of the WKB method.

## Disordered Chains

1. Consider resonant levels, consecutively coupled by matrix elements randomly distributed according to Poisson statistics, by analogy to the onedimensional elastic chain of randomly distributed masses (F. J. Dyson [68] The dynamics of a disordered linear chain (1953)) shown in Fig. 2.21 and find the mean density of states.
2. Consider the mean Green functions for the disordered chains (P. Lloyd[131] Exactly solvable model of electronic states in a three-dimensional disordered Hamiltonian: non-existence of localized states (1969)).
3. Consider quantum localization with the help of the scaling idea applied to transparencies and reflectivities of groups of consecutive levels in the disordered chain (P. W. Anderson, D. J. Thouless, E. Abrahams, and D. S. Fisher[8] New method for a scaling theory of localization (1980)).
4. For the same problem consider the role of phase relations in the scaling (L. I. Deych, A. A. Lisyansky, and B. L. Altshuler[59] Single parameter scaling in one-dimensional localization revisited (2000)).
5. Consider the density of states fora non-Hermitian chain (C. Mudry, P. W. Brouwer, B. I. Halperin, V. Gurarie, A. Zee[142] Density of states in the non-Hermitian Lloyd model (1998)).
6. Consider the effect of a time-dependent perturbation on the localized electron states in a random one-dimensional lattice (N.F. Mott and W. D. Twose[139] The theory of impurity conduction §7 (1961)).

Composite Systems. Relay of Bands

1. Consider the possibility of representating generic, multilevel systems as a level-band system and as a relay of isolated levels (E. Kyrölä, and M. Lindberg[122] Serial and parallel multilevel systems (1987)).
2. Represent the Shcrödinger equation describing the Hall effect (B. Huckestein[97] Scaling theory of the integer quantum Hall effect (1995)) as a relay of degenerate bands of Landau levels.
3. Derive the master equation for randomly coupled, multi-band systems (Sect. 8.1.2) with the help of a standard Feynman diagram approach (B. Carmeli, A. Nitzan[39] Kinetic equation for collisionless multiphoton excitation of large molecules (1978) and[40] Random coupling model for intramolecular dynamics (1980)).
4. Consider the limit of applicability of the master equation (D. L. Shepelyansky[164] Localization of diffusive excitation in multi-level systems (1987)).
5. Consider the one-dimensional chain of levels with time-dependent level positions (J. Karczmarek, M. Scott, and M. Ivanov[109] Two-color control of localization: From lattices to spin systems (1999)) in the quasienergy representation, as a relay of bands.
6. Describe a quantum dot in a microwave field (D. M. Basko, M. A. Skvortsov, and V. E. Kravtsov[13] Dynamic localization in quantum dots: analytical theory (2003)) as a relay of quasienergy bands.
7. Describe an electron in a nanotube (L.S. Levitov and A. M. Tsvelik[127] Narrow-gap Luttinger liquid in carbon nanotubes (2003)) in terms of a multiband system with four levels in each band.
8. Describe Arnold diffusion (V. Ya. Demikhovskii, F. M. Izrailev, and A. I. Malyshev[55] Manifestation of Arnold diffusion in quantum systems (2002)) as a process in a multiband system.
9. Consider multiphoton excitation of atoms above the ionization limit (N. B. Delone and M. V. Fedorov [52] Above-Threshold Ionization (1989), and [53] Multiphoton ionization of atoms: New effects (1989)).

## Effect of Correlations and Relay with Decay

1. Consider a relay of isolated levels coupled to a continuum (R. S. Burkey and C. D. Cantrell[37] Multichannel excitation of the quasi-continuum (1985)) as a multiband system with correlated matrix elements.
2. Consider the above-threshold ionization of atoms (Z. Deng and J. H. Eberly[58] Multiphoton absorption above ionization threshold by atoms in strong laser field (1985)) as a process in a multiband relay with completly correlated matrix elements.
3. Consider a Rydberg atom in a strong laser field (M. V. Fedorov and A. M. Movsesian[78] Interference suppression of photoionization of Rydberg atoms in a strong electromagnetic field (1988) and [79] Field-induced effects of narrowing of photoelectron spectra and stabilization of Rydberg atoms (1988)) as continuum-continuum transitions with a rank-1 interaction.
4. Consider a hydrogen atom in an intense RF field (G. Casati, B. V. Chirikov, and D. L. Shepelyansky[41] Quantum limitation for chaotic excitation of the hydrogen atom in a monochromatic field (1984)) as a relay of bands with correlated matrix elements. Estimate the typical time for ballistic propagation for this model.

## Localization in Relay of Bands, Field-Theory Methods

1. Determine the quantum localization border for Rydberg atoms in an RF field (G. Casati, B. V. Chirikov, I. Guarneri, and D. L. Shepelyansky[42] Dynamical stability of quantum "chaotic" motion in a hydrogen atom (1986)) and find the parameter gouverning the localization (G. Casati, I. Guarneri, F. Izrailev, and R. Scharf[43] Scaling behavior of localization in quantum chaos (1990)). Express this in terms of band random matrices (G. Casati, L. Molinari, and F. Izrailev[44] Scaling properties of band random matrices (1990)) and compare it to molecular models (G.Grahm and D. Shepelyansky[90] A solid-state model for photonic localization in molecular quasicontinua (1991)).
2. Consider compex systems with the help of the $\sigma$-model technique and the saddle-point approximation. (K. Efetov[74] Supersymmetry in disorder and chaos (1996), Chap. 4).
3. Consider random band matrices in terms of a $\sigma$-model (Y. V. Fyodorov and A. D. Mirlin[82] Scaling properties of localization in random band matrices: A $\sigma$-model approach (1991). Consider the transition to localization in such systems (A. D. Mirlin, Y. V. Fyodorov, F-M. Dittes J. Quezada, and
T. H. Seligman[138] Transition from localized to extanded eigenstates in the ensembles of power-law random banded matrices (1996)).
4. Consider the role of particle collisions on the localization (B. L. Altshuler, A. G. Aronov and D. E. Khmelnitsky[6] Effects of electron-electron collisions with small energy transfers on quantum localisation (1982)). Formulate the problem of binary collisions in terms of random matrices (D. L. Shepelyansky[165] Coherent propagation of two interacting particles in a random potential (1994)). Consider this problem with the help of a $\sigma$-model (Y. V. Fyodorov, and A. Mirlin[83] Statistical properties of random banded matrices with strongly fluctuating diagonal elements (1995)).
5. Consider the transition to localization for sparse band matrices ( Y. V. Fyodorov, A. D. Mirlin, and H-J. Sommers[85] A novel field theoretical approach to the Anderson localization: sparse random hopping model (1992)).
6. Consider the dynamics of elementary excitations in cold Rydberg atoms as a level coupled to a band of random walks (I. Mourachko et al.[140] Manybody effects in a frozen rydberg gas (1998)).

## Transformation of Complex Spectra

1. Consider the transformation of a non-degenerate spectrum perturbed by a small random matrix (P.Pechukas[149] Distribution of energy eigenvalues in the irregular spectrum (1983)). Derive the diffusion equation describing this process.
2. Derive a functional equation describing the transformation of a general multilevel system perturbed by a random matrix. (L.A. Pastur[147] On the spectrum of random matrices (1972)).
3. Derive the Pastur equation via a direct average over a random matrix perturbation of a deterministic quantum system (E.Brezin, S. Hikami, and A. Zee[33] Universal correlations for deterministic plus random Hamiltonians (1995)), and find the correlations of the level positions. Consider the corresponding time evolution (E. Brezin and S. Hikami [29] Spectral form factor in a random matrix theory (1997)). Consider the general case of level correlations (E. Brezin and S. Hakami[30] Extension of level-spacing universality (1997).
4. Trace the transition from the Breit-Wigner density of states to a Wigner semicircular profile for sparse band random matrices (Y. V. Fyodorov, O. A. Chubykalo, F. M. Izrailev and G. Casati[84] Wigner random banded matrices with sparse steucture: local spectral density of states (1996)).

5 . Derive a differential equation of motion for the mean density of states considering the mean square coupling of a random perturbation as a variable (F. Dyson[70] A class of matrix ensembles (1972)).
6. Consider the behavior near the gap closing (E.Brezin, S. Hikami[31] Universal singularity at the closure of a gap in a random matrix theory (1998)).
7. Find the transformation under perturbation by a random matrix of the following spectral profiles: (a) Density of states of a one-dimensional free particle $g(E)=2 \pi\left(1-E^{2}\right)^{-1 / 2}$; (b) Hubbard random alloys; (c) twodimensional free particle $g(E)=\pi \Theta(E)$; (d) a particle in the conduction band of a semiconductor with the kinetic energy $T(p)=\cos p$; (e) the Rydberg spectrum $E=1 / n^{2}$ for $n=2,3,4$ as given by the solution of the fourthorder algebraic equation.

## References

1. G. S. Agarwal, Control of decoherence and relaxation by frequency modulation of a heat bath, Phys. Rev. A 61, 013809 (2000)
2. V.M. Akulin, Contribution of weak and multiphoton transitions to the excitation of polyatomic molecules in an intense laser field, Sov. Phys. JETP 56, 533-538, (1982)
3. V.M. Akulin, Singularities of collisional vibrational relaxation processes in strongly excited spectrally complex polyatomic molecules, Sov. Phys. JETP 57, 774-780, (1983)
4. V. M. Akulin and N. V. Karlov, Intense resonant interactions in quantum electronics, Springer (1991)
5. L. Allen, J. H. Eberly, Optical resonance and two-level atoms, John Wiley and Sons (1975)
6. B. L. Altshuler, A. G. Aronov and D. E. Khmelnitsky, Effects of electronelectron collisions with small energy transfers on quantum localisation, J. Phys. C: Solid State Phys. 15, 7367-7386 (1982)
7. P. W. Anderson, Absence of diffusion in certain random lattices, Phys. Rev. 109, 1492-1505 (1957)
8. P. W. Anderson, D. J. Thouless, E. Abrahams, and D. S. Fisher New method for a scaling theory of localization, Phys. Rev. B 22, 3519-3526 (1980)
9. V. I. Arnold, A. Weinstein, and K. Vogtmann, Mathematical methods of classical mechanics, Springer Verlag (1989)
10. G. A. Askaryan, Exitation and dissociation of molecules in intense light field, Sov. Phys. JETP 19, 273-274 (1964)
11. I. Sh. Averbukh and N. F. Perelman, Fractional revivals: universality in the long-term evolution of quantum wave packets beyong the correspondence principle dynamics, Physics Letters A ,139, 449-454 (1989)
12. A. Barut and R. Raczka, Theory of group representations and applications, Polish Scientific Publisher, Warszava (1977)
13. D. M. Basko, M. A. Skvortsov, and V. E. Kravtsov, Dynamic localization in quantum dots: analytical theory, Phys. Rev. Lett. 90, 096801-096804 (2003)
14. C. W. J. Beenakker, Random matrix theory of quantum transport, Rev. Mod. Phys 69, 731-808 (1997)
15. V. L. Berezinskii, Kinetics of a quantum particle in a one-dimensional random potential, Sov. Phys. JETP 38, 620-627 (1974)
16. M. V. Bery and M. Tabor, Poissonian statistics and clustering, Proc. R. Soc. Lond. 356, 375-394 (1977)
17. M. Bixon and J. Jortner, Intramolecular radiationless transitions, J. Chem. Phys. 48, 715-726 (1968)
18. M. Bixon and J. Jortner, Long radiative lifetimes of small molecules, J. Chem. Phys. 50, 3284-3290 (1969)
19. Ya. M. Blanter, A. D. Mirlin, and B. A. Muzykantskii, Quantum billiards with surface scattering: Ballistic sigma-model approach, Phys. Rev. Lett. 80, 4161 (1998)
20. F. Bloch, Nuclear induction, Phys. Rev. 70, 460 (1946)
21. F. Bloch, Generalized theory of relaxation, Phys. Rev. 105, 1206-122 (1957)
22. F. Bloch and I. I. Rabi, Atoms in variable magnetic fields, Rev. Mod. Phys. 17, 237-244 (1945)
23. N. Blombergen, Nonlinerar optics, Benjamin (1965)
24. O. Bohigas M. J. Giannoni, and C. Schmit, Characterization of chaotic quantum spectra and universality of level fluctuation laws, Phys. Rev. Lett. 52, 1-4 (1984)
25. E. B. Bogomolny, U. Gerland, and C. Schmit, Models of intermediate spectral statistics, Phys. Rev. E 59, R1315-R1318 (1999)
26. E. B. Bogomolny, U. Gerland, and C. Schmit, Singular statistics, Phys. Rev. E 63, 036206 (2001)
27. R. E. Borland, The nature of the electronic states in disordered onedimensional systems, Proc. Roy. Soc. Lond. A 274, 529-545 (1963)
28. G. Breit and E. Wigner, Capture of slow neutrons, Phys. Rev. 40, 519-531 (1936)
29. E.Brezin, S. Hikami, Spectral form factor in a random matrix theory, Phys. Rev. E 55, 4067-4083 (1997)
30. E. Brezin and S. Hakami, Extension of level-spacing universality, Phys. Rev. E 56, 264-269 (1997)
31. E. Brezin, S. Hikami, Universal singularity at the closure of a gap in a random matrix theory, Phys. Rev. E, 57, 4140-4149 (1998)
32. E. Brezin and S. Hikami, New correlation functions for random matrices and integrals over supergroups, J. Physics A 36, 711-751 (2003)
33. E.Brezin, S. Hikami, and A.Zee, Universal correlations for the deterministic plus random Hamiltonians, Physical Review E 51, 5442-5452 (1995)
34. P. W. Brumer, M. Shapiro, Principles of the quantum control of molecular processes, John Wiley and Sons, (2003)
35. F. V. Bunkin, A. M. Prokhorov, The excitation and ionization of atoms in a strong radiation field, Sov. Phys. JETP 19, 739-743 (1964)
36. F. V. Bunkin, A. M. Prokhorov, and R. V. Karapetyan, Dissociation of molecules in a strong radiation field, Sov. Phys. JETP 20, 145-148 (1965)
37. R. S. Burkey and C. D. Cantrell, Multichannel excitation of the quasicontinuum, J. Opt. Soc. Am. B 2, 451-457 (1985)
38. A. G. Butkovskiy and Yu. I. Samoilenko, Control of quantum-mechanical processes and systems, Kluwer Academic Publishers, (1990)
39. B. Carmeli, A. Nitzan, Kinetic equation for collisionless multiphoton excitation of large molecules, Chem. Phys. Lett. 62, 457-460 (1978)
40. B. Carmeli, A. Nitzan, Random coupling model for intramolecular dynamics, J. Chem. Phys. 72, 2054-2080 (1980)
41. G. Casati, B. V. Chiricov, and D. L. Shepelyansky, Quantum limitation for chaotic excitation of the hydrogen atom in a monochromatic field, Phys. Rev. Lett. 53 2525-2528, (1984)
42. G. Casati, B. V. Chirikov, I. Guarneri, and D. L. Shepelyansky, Dynamical stability of quantum chaotic motion in a hydrogen atom, Phys. Rev. Lett. 23, 2437-2440 (1986)
43. G. Casati, I. Guarneri, F. Izrailev, and R. Scharf, Scaling behavior of localization in quantum chaos, Phys. Rev. Lett. 64, 5-8, (1990)
44. G. Casati, L. Molinari, and F. Izrailev, Scaling properties of band random matrices, Phys. Rev. Lett. 64, 1851-1854 (1990)
45. J. T. Chalker, I. V. Lerner and R. A. Smith, Random walks through the ensemble: Linking spectral statistics with wave-function correlations in disordered metals, Phys. Rev. Lett. 77, 554-557 (1996)
46. B. V. Chirikov, Resonance processes in magnetic traps, J. Nucl. Energy, Part C Plasma Phys. 1, 253-260 (1960)
47. B.V.Chirikov, D.L.Shepelyansky, Correlation properties of dynamical chaos in hamiltonian systems, Physica D 13, 395-400 (1984)
48. J. I. Cirac and P. Zoller, Quantum computation with cold trapped ions, Phys. Rev. Lett., 74, 4091-4094, (1995).
49. C. Cohen-Tanoudji, B. Diu, and F. Laloe, Quantum mechanics, Wiley, (1977)
50. C. Cohen-Tannoudji, J. Dupont-Roc, and G. Grynberg, Atom-photon interactions : Basic processes and applications, Wiley-Interscience, (1998)
51. D. Delande and J. Zakrzewski, Spontaneous emission of nondispersive Rydberg wave packets, Phys. Rev. A 58, 446-477 (1998)
52. N. B. Delone and M. V. Fedorov, Above-threshold ionization, Progr. Quant. Electron. 13, 267-298 (1989)
53. N. B. Delone and M. V. Fedorov, New effects in the multiphoton ionization of atoms, Sov. Phys. Uspekhi 32, 500-520 (1989)
54. N.B. Delone, V.P. Krainov, Atoms in strong light fields, Springer Verlag; (1985)
55. V. Ya. Demikhovskii, F. M. Izrailev, and A. I. Malyshev, Manifestation of Arnold diffusion in quantum systems, Phys. Rev. Lett. 88, 154101 (2002)
56. Yu. N. Demkov, Charge transfer at small resonance defects, JETP 18, 138-142 (1964)
57. Yu. N. Demkov and V. I. Osherov, Stationary and nonstationary problems in quantum mechanics that can be solved by means of contour integration, Sov. Phys. JETP, 26, 916-921 (1968)
58. Z. Deng and J. H. Eberly, Multiphoton absorption above ionization threshold by atoms in strong laser field, J. Opt. Soc. Am. B 2, 486-493 (1985)
59. L. I. Deych, A. A. Lisyansky, and B. L. Altshuler, Single parameter scaling in one-dimensional localization revisited, Phys. Rev. Lett. 84, 2678-2681, (2000)
60. R. H. Dicke, Coherence in spontaneous radiation processes, Phys. Rev. 93, 99 (1954)
61. V. V. Dodonov, Universal integrals of motion and universal invariants of quantum systems, J. Phys. A 33, 7721-7738 (2000)
62. V. V. Dodonov and V. I. Man'ko, Phase space eigenfunctions of multidimensional quadratic hamiltonians, Physica A 137, 306-316 (1986)
63. V. V. Dodonov, V. I. Manko, and D. E. Nikonov, Exact propagators for timedependent Coulomb, delta and other potentials, Phys. Lett.A 162, 359-364 (1992)
64. G. Drobny, I. Jex, and V.Buzek, Mode entanglement in nondegenerate downconversion with quantized pump, Phys. Rev. A 48, 569-579 (1993)
65. Yu. V. Dubrovskii, M. V. Fedorov, and M. Yu. Ivanov, Rydberg atom ionization by an intense laser pulse with smooth time envelope: a model of two closely spaced levels coupled to continuum, Laser Physics, 1992, v. 2, pp. 288-298, (1992)
66. Ph. Durand, I. Paladrova and F. X. Cadea, Theory of Fano profiles, J. Phys. B 34, 1953-1986 (2001)
67. A. M. Dykhne, Adiabatic perturbation of the descrete spectrum states, Sov. Phys. JETP 14, 941-943 (1962)
68. F. J. Dyson, The dynamics of a disordered linear chain, Phys. Rev. 92, 13311338 (1953)
69. F. Dyson, Statistical theory of the energy levels of complex sustems, J. Math. Phys. 3, 140-175 (1962)
70. F. Dyson, A class of matrix ensembles, J. Math. Phys. 13, 90-97 (1972)
71. J. Eberly, N. B. Narozhny, and J. J. Sanchez-Mondregon, Periodic sopntaneous collapse and revivals in a simple quantum model, Phys. Rev. Lett. 44, 13231325 (1980)
72. B. Eckhardt, Quantum mechanics of classically non-integrable systems, Phys. Rep. 163, 205-297 (1988)
73. A. Einstein, Zur Quantentheorie der Strahlung, Physikalische Zeitschrift 18 (1917), 121-128; The Collected Papers of Albert Einstein, Volume 6: The Berlin Years: Writings, 1914-1917. (English translation supplement) University Press, Princeton
74. K. Efetov, Supersymmetry in disorder and chaos, Cambridge University Press, (1997)
75. U. Fano, Sullo spettro di assorbimento dei gas nobili presso il limite dello spettro d'arco, Nuovo Cimento, N. S., 12, 154-161 (1935).
76. U. Fano, Effects of configuration interaction on intensities and phase shifts, Phys. Rev. 124, 1866-1878 (1961)
77. M. V. Fedorov, M. Yu. Ivanov and A. M. Movsesian, Strong-field photoionization of an initially excited hydrogen atom: Formation of Rydberg wavepacket, its structure and trapping of population at Rydberg levels, J. Phys. B 23, 2245S2257S (1990)
78. M. V. Fedorov and A. M. Movsesian, Interference suppression of photoionization of Rydberg atoms in strong electromagnetic field, J. Opt. Soc. Am. B 6, 928-936 (1988)
79. M. V. Fedorov and A. M. Movsesian, Field-induced effects of narrowing of photoelectron spectra and stabilization of Rydberg atoms, J. Phys. B 21, L155 (1988)
80. E. Fermi, Quantum theory of radiation, Rev. Mod. Phys. 4, 87-133 (1932)
81. G. Floquet, Sur les équations différentialles linéaires à coefficients périodiques, Ann. de l'Ecole Normale 2-e serie, 12 47-88 (1889)
82. Y. V. Fyodorov and A. D. Mirlin, Scaling properties of localization in random band matrices: A $\sigma$-model approach, Phys. Rev. Lett. 67, 2405-2409 (1991)
83. Y. V. Fyodorov, and A. Mirlin, Statistical properties of random banded matrices with strongly fluctuating diagonal elements, Phys. Rev. B 52, R11580R11583 (1995)
84. Y. V. Fyodorov, O. A. Chubykalo, F. M. Izrailev and G. Casati, Wigner random banded matrices with sparse steucture: local spectral density of states, Phys. Rev. Lett. 76, 1603-1606 (1996)
85. Y. V. Fyodorov, A. D. Mirlin, and H-J. Sommers, A novel field theoretical approach to the Anderson localization : sparse random hopping model, J. Phys. I France 2, 1571-1605 (1992)
86. R. G. Gallager, Information theory and reliable communications, John Wiley and sons, Inc. (1968).
87. T. F. Gallagher, Rydberg atoms, Cambridge University Press; (1994)
88. M. Gavrila and J. Z. Kaminski, Free-free transitions in intense high-frequency laser fields, Phys. Rev. Lett. 52, 613-616 (1984)
89. P. K. Ghosh, Ion traps, Clarendon Pr; (1995)
90. G.Grahm and D. Shepelyansky, A solid-state model for photonic localization in molecular quasi-continua, Europhys. Lett. 14, 211-215 (1991)
91. M. Gruebele, Intramolecular vibrational dephasing obays a power law at intermediate times, Proc. Natl. Acad. Sci. USA 95, 5965-5970, (1998)
92. S. Guérin and H. R. Jauslin, Control of quantum systems by pulses: Adiabatic Floquet theory, Adv. Chem. Phys. 125,1-75, (2003).
93. M. Gutzwiller, Chaos in classical and quantum mechanics, Springer Verlag (1990).
94. F. Haake, Quantum signatures of chaos, Springer Verlag, (2001)
95. E. L. Hahn, Spin echos, Phys. Rev. 80, 580-594 (1950)
96. J. Hubbard, Electron correlations in narrow energy bands, Proc. R. Soc. London, Ser. A 276 238-257(1963); 277, 237-259 (1964); 281, 401-419 (1964)
97. B. Huckestein, Scaling theory of the integer quantum Hall effect, Rev. Mod. Phys. 67, 375-396 (1995)
98. K. Husimi, Miscellanea in elementary Quantum Mechanics, Prog. Theor. Phys. 9, 238-244; 381-402 (1953)
99. J. Imry, Introduction to mesoscopic physics, Oxford University Press (2002)
100. K. Ishii, Localization of eigenstates and transport phenomena in the onedimensional disordered systems, Supl. Prog. Theor. Phys. 53, 77-138, (1972)
101. M. Yu. Ivanov, Suppression of resonant multiphoton ionization via Rydberg states, Phys. Rev. A 49, 1165-1170 (1994)
102. F. M. Izrailev, Simple models of quantum chaos: Spectrum and eigenfunctions, Physics Reports 5-6, 299-392, (1990)
103. D. Jaksch, H. Briegel, J. I. Cirac, C. W. Gardiner, and P. Zoller, Entanglement of atoms via cold controlled collisions, Phys. Rev. Lett. 82, 1975-1978 (1999)
104. J. Javanainen and E. Kyrölä, Long-time limit of a quasicontinuum model, Opt. Comm. 56, 17-21 (1985)
105. J. D. Joannopoulos, R. D. Meade, and J. N. Winn, Photonic crystals, Princeton Univ Pr; (1995).
106. Journal of Physics A 36, Number 12, 28 March (2003)
107. Ch. Jungen, Molecular applications of quantum defect theory, Institute of Physics Pub, (1996)
108. V. Jurdevic and H. J. Sussmann, Control systems on Lie groups, Journal of Differential Equations, 12, 313-329 (1972)
109. J. Karczmarek, M. Scott, and M. Ivanov, Two-color control of localization: From lattices to spin systems, Phys. Rev. A 60, R4225-R4228 (1999)
110. R. Karplus and J. Schwinger, A note on saturation in microwave spectroscopy, Phys. Rev. 73, 1020-1026 (1948)
111. P. L. Knight, M. A. Lauder and B. J. Dalton, Laser-induced continuum structure, Phys. Reports 190, 1-61 (1990)
112. A. G. Kofman, Theory of single-photon bound-free transitions: Extension of the pole approximation, J. Phys. B 30, 5141-5156 (1997)
113. A. G. Kofman, Relaxation of a two-level system strongly coupled to a reservoir: Anomalously slow decoherence, Phys. Rev. A 64, 033809 (2001)
114. A. G. Kofman and G. Kurizki, Universal dynamical control of quantum mechanical decay: Modulation of the coupling to the continuum, Phys. Rev. Lett. 87, 270405 (2001)
115. V. A. Kostelecky, V. I. Manko, M. M. Nieto, and D. R. Traux, Supersymmetry and time-dependent Landau system, Phys. Rev. A 48, 951-963 (1993)
116. V. E. Kravtsov, Time reversal symmetry breaking by ac field: Effect of commensurability in frequency domain, Pramana - journal of physics 58, 183-193 (2002)
117. V. E. Kravtsov, I. V. Lerner, B. L. Altshuler and A. G. Aronov, Universal spectral correlations at the mobility edge, Phys. Rev. Lett. 72, 888-891 (1994)
118. R. Kubo, A general expression for the conductivity tensor, Canadian J. Phys. 34,1274-1277 (1956)
119. G. Kurizki, A. G. Kofman, and V. Yudson, Resonant photon exchange by atom pairs in high-Q cavities, Phys. Rev. A 53, R35-R38 (1996)
120. D. Kuznetsov, A. Bulgak and G. D. Dang, Quantum Lévy process and fractional kinetics, Phys. Rev. Lett. 82, 1136-1139 (1999)
121. E. Kyrölä and J. Eberly, Quasicontinuum effects in molecular exciation, J. Chem. Phys. 82, 1841-1853 (1985)
122. E. Kyrölä and M. Lindberg, Serial and parallel multilevel systems, Phys. Rev. A 35, 4207-4225 (1987)
123. L. D.Landau, Das Dämpfungsproblem in der Wellenmechanik, Zeitschrift für Physik, Bd XLV, 430-441 (1927)
124. L. D. Landau and E. M. Lifshits, Quantum mechanics, ButterworthHeinemann (1975-1999)
125. C. Lechtle, I. Sh. Averbukh, and W. P. Schleich, Generic structure of multilevel quantum beats, Phys. Rev. Lett. 77, 3999-4002 (1996)
126. R. D. Levine, Representation of one-dimensional motion in a Morse potential by a quadratic hamiltonian, Chem. Phys. Lett. 95, 87-90 (1983)
127. L.S. Levitov and A. M. Tsvelik, Narrow-gap Luttinger liquid in carbon nanotubes, Phys. Rev. Lett. 90, 016401 (2003)
128. A.J. Lichtenberg, M.A. Liberman, Regular and stochastic motion, Springer Verlag, (1982)
129. I.M. Lifshits, S.A. Gredeskul and L.A. Pastur, Introduction to the theory of disordered systems, John Wiley and Sons (1988)
130. M. Lindberg and E. Kyrölä, Infinite multilevel system coupled to a quantumfield mode, J. Opt. Soc. Am. B, 9, 595-604 (1992)
131. P. Lloyd, Exactly solvable model of electronic states in a three-dimensional disordered Hamiltonian: non-existence of localized states, J. Phys. C: Solid State Phys. 2, 1717-1725 (1969)
132. A.Makarov, Coherent excitation of equidistant multilevel systems in a resonant monochromatic field, Sov. Phys. JETP 45, 918-924 (1977)
133. H.S.W. Massey and N. F. Mott, The theory of atomic collisions, Clarendon Pr (1997)
134. J. Mathews and R. L. Walker, Mathematical methods of physics, AddisonWesley (1971)
135. S. L. McCall and E. L. Hahn, Self-induced transparency, Phys. Rev. 183, 457485 (1969)
136. M. L. Metha, Random matrices and statistical theory of energy levels, Academic Press, New York and London, (1967)
137. P. Meystre, Atom optics, Springer Verlag (2001)
138. A. D. Mirlin, Y. V. Fyodorov, F-M. Dittes J. Quezada, and T. H. Seligman, Transition from localized to extanded eigenstates in the ensembles of power-law random banded matrices, Phys.Rev. E 54, 3221-3230 (1996)
139. N.F. Mott and W. D. Twose, The theory of impurity conduction, Adv. Phys. 10, 107-163 (1961)
140. I. Mourachko, D. Comparat, F. de Tomasi, A. Fioretti, P. Nosbaum, V.M. Akulin, and P. Pillet, Many-body effects in a frozen Rydberg gas, Phys.Rev. Lett. 80, 253-256 (1998)
141. L. Mower, Decay theory of closely coupled unstable states, Phys. Rev. 142, 799-816 (1966)
142. C. Mudry, P. W. Brouwer, B. I. Halperin, V. Gurarie, A. Zee, Density of states in the non-Hermitian Lloyd model, Phys. Rev. B 58, 13539 (1998)
143. S. Mukamel and J. Jortner, A model for isotop separation via molecular multiphoton dissociation, Chem. Phys. Lett., 150-156 (1976)
144. R. Narevich, R. E. Prange, O. Zaitsev, Square billiard with an Aharonov-Bohm flux line, Physica E 9, 578-582 (2001)
145. E. E. Nikitin, Theory elementary atomic and molecular processes, Oxford University Press, (1985)
146. E. E. Nikitin S. Ya. Umanskii, Theory of slow atomic collisions, Springer Verlag, (1984)
147. L.A. Pastur, On the spectrum of random matrices, Theor. Math. Phys. (USSR) 10, 67-74 (1972)
148. P. Pechukas, Quantum chaos in the irregular spectrum, Chem. Phys. Letters 86, 553-557 (1982)
149. P. Pechukas, Distribution of energy eigenvalues in the irregular spectrum, Phys. Rev. Lett. 51, 943-946 (1983)
150. A. M. Perelomov, Integrable systems of classical mechanics and Lie algebras, Springer Verlag; (1990); Generalized coherent states and their applications, Springer Verlag; (1986)
151. C. E. Porter, Statistical theories of spectra: fluctuations, Academic Press, New York and London, (1965)
152. C.E. Porter and R.G. Thomas, Fluctuation of nuclear reaction width, Phys. Rev. 104, 483-491 (1956)
153. M. Quack, Theory of unimolecular reactions induced by monochromatic infrared radiation, J. Chem. Phys. 69 1282-1307 (1978)
154. I. I. Rabi, N. F. Ramsey, and J. Schwinger, Use of rotating coordinates in magnetic resonance problem, Rev. Mod. Phys. 26, 167-171 (1954)
155. L. E. Reichl, The transition to chaos, Springer Verlag (1992)
156. C. K. Rhodes, A. Szöke, and A. Javan, The influence of level degeneracy on the self-induced transparency effect, Phys. Rev. Lett. 21, 1151-1155 (1968)
157. V. I. Ritus, Shift and splitting of atomic energy levels by the field of an electromagnetic wave, Sov. Phys. JETP 24, 1041-1044 (1967)
158. M. Sargent, W. E. Lamb, M. O. Scully, Laser physics, Addison-Wesley Publishing; (1987)
159. K. J. Schafer and K. C. Kulander, Phase dependent effects in multiphoton ionization induced by a laser field and its second harmonic, Phys. Rev. A 45, 8026-8033 (1992)
160. E. W. Schlag, ZEKE Spectroscopy, Cambridge University Press (1998)
161. W.Schleich, Quantum optics in phase space, John Wiley (2001)
162. S. A. Schofield and P. G. Wolynes, A scaling perspective of quantum energy flow in molecules, J. Chem. Phys. 98, 1123-1131 (1992)
163. P. Seba, Wave chaos in singular quantum billiard, Phys. Rev. Lett. 64, 18551858 (1990)
164. D. L. Shepelyansky, Localization of diffusive excitation in multi-level systems, Physica D: Nonlinear Phenomena, 28, 103-114 (1987)
165. D. L. Shepelyansky, Coherent propagation of two interacting particles in a random potential, Phys. Rev. Lett. 73, 2607-2610 (1994)
166. B. W. Shore and J. H. Eberly, Analytic approximations in multi-level excitation theory, Opt. Comm. 24, 83-88 (1978)
167. Ya. G. Sinai, On the concept of entropy for a dynamic system, Sov. Phys. Doklady, 124, 768-771 (1959)
168. Ya. G. Sinai, Introduction to ergodic theory, Princeton Univ Pr. (1976)
169. Ya. G. Sinai, Anderson localization for one-dimensional difference Schrödinger operator with quasi-periodic potential, J. Stat. Phys. 46, 861-909 (1987)
170. Ya. G. Sinai, Topics in ergodic theory, Princeton Univ Pr, (1994)
171. Ya. G. Sinai and A. B. Soshnikov, A refinement of Wigner's semicircle in a neighbouhood of the spectrum edge for random symmetric matrices, Functional Analysis and Applications 32, 114-131 (1998)
172. I. I. Sobelman, L. A. Vainshtain, and E. A. Yukov, Excitation of atoms and broadening of spectral lines, Springer Verlag (1981)
173. S. Stenholm, Quantum theory of electromagnetic fields interacting with atoms and molecules, Physics Reports, 6, 1-122 (1973)
174. D. J. Thouless, Maximum metallic resistance in thin wires, Phys. Rev. Lett 39, 1167-1169 (1977)
175. V. Weisskopf und E. Wigner, Berechnung der natürlichen Linienbreite auf Grund der Diracschen Lichttheorie, Z. Phys. 63, 54-73 (1930)
176. E. Wigner, On a class of analytic functions from the quantum theory of collisions, Ann. Math. 53; 36 (1951)
177. E. Wigner, On the distribution of the roors of certain symmetric matrices, Ann. Math 65, 203-207 (1957)
178. V. A. Yurovsky, A. Ben-Reuven, Saturated transition in exactly soluble models of two-state curve crossing with time-dependent potentials, Phys. Rev. A 60, 4561-4566 (1999)
179. V. A. Yurovsky, A. Ben-Reuven, P. S. Julienne, and Y. B. Brand, Conterintuitive transitions in multistate curve crossing involving linear potentials, J. Phys. B 32, 1845-1857 (1999)
180. G. M. Zaslavski, R. Z. Sagdeev, D. A. Usikov, and A. A. Chernikov, Weak chaos and quasi-regular patterns, Cambridge University Press (1991)
181. Ya. B. Zeldovich, The quasienergy of a quantum mechanical system subjected to a periodic action, Sov. Phys. JETP 24, 1006-1008 (1967)
182. Ya. B. Zeldovich, Scattering and emission of a quantum system in a strong electromagnetic wave, Sov. Phys. Uspekhi 16, 427-433 (1973)
183. M. Zirnbauer, Anderson localization and non linear sigma-model with graded symmetry, Nuclear Physics B 265, 375-408 (1986)

## Index

Adiabatic (see also WKB)
and diabatic terms, 197
Dykhne formula, 217
frame, 216
non-adiabatic interaction, 36
non-adiabatic transition, 216
separation, 21, 41
Algebraic
condition of controllability, 285, 294, 296
properties, 187, 188
symmetry and solubility, 195, 311
Ancilla, 302
Anharmonicity
causing fractional revivals, 90, 93, 97
of molecular vibrations, 22
Arnold diffusion, 10
Baker-Campbell-Hausdorff formula, 188, 191, 220, 294
Bloch vector, 85, 303
Bohigas-Giannoni-Schmit conjecture, 7, 9
Born-Oppenheimer parameter, 35, 41
Brackets generation condition, 294
Breit-Wigner distribution, 15
Bright state, 170
Cantor set, 108
Causality principle and analytical properties, 181
Clusters, 18
metallic, 43, 46
molecular, 41
Coherent damping, 147, 236
Complexity
criterion, 11
Compound system, 30, 83
and environment, 83
and measurement, 83
Control, 284
algebraic conditions of, 285
bang-bang, 288
by a periodic field, 293
goals, 284
of coherence loss, 302
of non-holonom systems, 299
of quantum state, 288
of unitary evolution, 290
partial, 296
Cooperative
behavior, 76
coupling, 76
number of levels in resonance, 76
Correlations, 161
and rank 1,161
and rank of interaction, 123, 159
and stabilization, 178
between coupling and energy, 124, 179
complete, 165
Dark state, 170
Decay
bi-exponential, 166
exponential, 73
in Landau-Zener problem, 199
incomplete, $62,114,392,412$
non-exponential, 61, 106, 343, 408, 413
power law, 113, 116, 121, 186, 411
rate, 73
to a band of random walks, 392
Demkov-Osherov problem, 13, 222
Density matrix, 84
equations for, 85

Density of states
averaged, 43
for fractal structure, 431, 435
randomly perturbed, 437
for small random multiplets, 430
Wigner distribution for, 149
Diabatic terms, 197
Diagrams
for amplitudes, 124, 133
for populations, 127, 137
of tree-like topology, 132
self-intersecting, 128
Dicke effect, 56
Dipole-dipole interaction, 29
Edges of bands, 61
and descrete levels, 74
and non-adiabatic transition, 227
double passage of, 231
Energy exchange
in atoms, 58
vibrational, 29
Ensemble
average, 2, 3, 403
averaged amplitudes for disordered relay, 333
averaged density of states, 46
for disordered relay, 335
circular
orthogonal, 5, 14
simplectic, 5
unitary, 5, 14
Gaussian but not Wigner, 124
Gaussian(Wigner), 149, 422
orthogonal, 4, 251
simplectic, 4
unitary, 4, 251
non-Wigner, 106
Poissonian, 5, 106, 108, 400, 403
Environment, 83, 302
Euler trails, 397
number, 424
Factorized
complex expressions by integral representation, 108
number of Euler trails, 426
structure of coupling, 169
structure of Hamiltonians, 47
via Gaussian and Grassmann integrals, 401
Fano
problem, 68, 78
profile, 61, 81
Fermi golden rule, 11, 68
Feynmann path integral
for inverse three diagonal determinants, 346
Field theory method
for disordered systems, 344
bosonic fields, 344
fermionic fields, 354, 355
Floquet
approach, 193
theorem, 11
Fourier-Laplace transformation, 63
Fractals, 53
density of states, 431,435
multi, 60
randomly perturbed, 437
Graded variables (see Grassmann), 426
Grassmann variables (nilpotent, graded), 357, 359, 400, 426

Hamming condition, 306
Heisenberg return time, 10, 61, 72, 75, $83,88,224,410,422$
Hierarchical structure, 15
Holonom
constraint, 295
non-holonom, 295
partially, 295
sufficient criterion on non-
holonomicity, 296
system, 294
Injection(suply) of population, 61, 64, 81, 88
Integration over a matrix, 366
K-degeneracy, 25
Kapiza-Dirac potential, 48
Kolmogorov-Arnold-Moser tori, 8
Landau-Zener
formula, 199
problem, 13, 187, 194, 214
transition
at the continuum edge, 227
to decaying states, 199
Landau-Zener
transition
in the presence of coherence loss, 202
Laplace contour integral
choice of contour, 223
method, 187, 205, 222
Level-band
continuouse, 72
coupled to a random walk continuum, 397
exponentially rising coupling, 232
inhomogeneous, 105
moving degenerate level, 238
moving level, 222, 227
time-dependent, 221
Lie algebras and groups, 6, 187, 188
$\operatorname{su}(1,1), 196,294,314$
su(2), 195, 294, 321, 323
Linear absorption
cross-sections, 48
Liouville equation
for relay-like disordered system, 352
Lloyd chain, 14
Localization, 14, 52, 332
and incomplete decay, 392
measure, 53
Long time asymptotic, 88
for disordered relay, 340
for Level-band system, 111
Lorentzian (Cauchy) profile, 15, 74, 142
Master equation, 86
Measurement, 83
continuous, 87
elementary act of, 86
model of meter, 83
spectral characteristics, 87
Zeno effect, 87
Multiphoton
excitation, 25
processes, 27
resonances, 17
Number of resonant levels, 75
and cooperative behavior, 76,114
expectation value, $106,112,116$
Poisson summation formula, 92, 102, 264
Polyatomic molecule, 22, 27, 29
Population distribution
amongst bands, 143, 147, 174
amongst bands of relay
in the presence of correlations, 391
and interaction time dependence, 265
bi-harmonic perturbation by random matrix, 270
for band and degenerate level, 151
harmonic perturbation by random matrix, 261
abrupt switching, 261
adiabatic switching, 263
in a disordered chain, 369
over band, 140, 142, 154, 172, 186
in the presence of correlations in relay, 392
perturbed by time dependent random matrix, 246, 250
over two bands, 282
residual, 114, 413
role of the quantum field statistics, 279
two bands in a periodic field, 277
Quasicontinuum
criterion, 11
Quasienergies, 193
Rabi
frequency, 66, 392
oscillations, 61, 66, 107, 208
Random walk
and incomplete decay, 392
and population amplitude returns, 395
interference of returns, 397
non-returnable, 407
returnable, 407, 408
$\delta$-correlated, 407, 411
finite correlation time, 412
returns and interference, 407
Random walks
and coherent behavior, 392
Rank
and correlations, 123
of interaction, 159, 170
Recurrences, 10, 68, 72, 88, 91, 93, 244, 392
and random walks, 395
effect on the state density profiles, 422
shape, 95
Relaxation, 83
elementary act of, 83,86
longitudinal, 203
spectral characteristics, 87
transversal, 203
Relay-like system
bands of levels, 373
correlated coupling of bands, 387
degenerate bands with random coupling, 376
isolated levels, 309
decreasing coupling, 324
general soluble case, 317
harmonic oscillator, 311
increasing coupling, 322
Raman and linear pumping of oscillator, 316
Raman pumping of oscillator, 313
smooth varation of the parameters, 325
with uniform coupling and linear detuning, 310
WKB approximation, 326
non degenerate bands with random coupling, 380
tunneling probability, 329
with disordered parameters, 331, 360
tunneling transparency, 349
WKB for non degenerate bands, 384
Renormalization
and topology, 133
for fractals, 431
of energies, 135, 140, 141
Repulsion of levels, 61
Resonant
approximation, 20, 277
Response
to bi-harmonic perturbation, 14, 273
to harmonic perturbation, 265
abrupt switching, 265
adiabatic switching, 267
to perturbation
by a random matrix, 250
to perturbation by a random matrix, 14
Revivals, 10, 61, 68, 88, 95, 244, 392
and Erenfest theorem, 104
and the classical limit, 104
complete, 98
distribution of amplitudes for
fractional, 101
fractional, 61, 68, 88, 99
in irregular spectrum, 98
partial, 98
shape, 98
wave packets of fractional, 101
wavepacket, 97
Rotating wave
approximation, 62
representation, 68
Rydberg
atom, 17, 18, 26, 58
gas, 58
molecule, 18, 38, 40
Schrödinger cat, 61, 100
Separability
of holonom systems, 295
of molecular motion, 23
Sinai billiard, 8, 45, 47
Squeezed states, 315
Stabilization
for correlated couplings, 178
Statistical approach, 20
Statistics
and Wigner semicircle, 149
of deviations, 3
of inter-level distances, 6
Poissonian, 400
Supply of probability amplitude
(injection), 63, 70, 74
Susceptibility
to perturbation
abrupt, 266
adiabatic, 269
harmonic, 265
Time dependent system, 190
bi-harmonic perturbation by random matrix, 270
Floquet approach, 190
harmonic perturbation by random matrix, 261
Landau-Zener, 194
with decay, 199
with dephasing, 202
perturbation proportional to a random matrix, 244
Time-dependent system
adiabatic, 213
and Dykhne formula, 213
complex, 237
controllable, 285
Demkov-Osherov, 222
level-band, 221
multilevel, 218, 222
quasienergy, 193
two levels coupled by a pulse, 208
two-level system, 194
two-levels and exponential coupling, 210
Time-ordering, 187
Transformation
of absorption lines, 420
of spectral state density, 413, 414
rule for random perturbation, 415
Tunneling transparency
for disordered relay of levels, 349
Two level-band systems, 176
Two-band system, 123
band and degenerate level, 151
correlated coupling, 169
degenerate, 145
in a periodic field, 277
nondegenerate, 139
Two-frequency perturbation (biharmonic)
by random matrix, 270
of a two-band system, 14
Two-level system, 62
control, 286
excitation by pulse, 208
repulsion of levels, 6,61
time dependent, 194
with exponentially rising coupling, 210

Universal features in time dependencies, 421

Wavepacket
and Erenfest theorem, 104
and fractional revivals, 101
and recurrences, 93
for relay-like system of isolated levels, 327
of revivals, 97,98
representation, 77
Wigner semicircle, 149
WKB approximation, 14
applicability for relay, 330
Dykhne formula, 213, 217
for adiabatic systems, 213
for relay of non degenerate bands, 384
for relay-like system of isolated levels, 326

Zeno effect, 87, 302
and coherence protection, 304

