Chapter 6. Perturbative Approaches

This chapter discusses various perturbative approaches to problems of quantum mechanics, and their simplest but important applications including the fine structure of atomic energy levels, and the effects of external dc and ac electric and magnetic fields of these levels. It continues with a discussion of the perturbation theory of transitions to continuous spectrum and the Golden Rule of quantum mechanics, which will naturally bring us to the issue of open quantum systems – to be discussed in the next chapter.

6.1. Eigenproblems

Unfortunately, only a few problems of quantum mechanics may be solved exactly in the analytical form. Actually, in the previous chapters we have solved a substantial part of such problems for a single particle, and for multiparticle problems, the exactly solvable cases are even more rare. However, most practical problems of physics feature a certain small parameter, and this smallness may be exploited by various approximate analytical methods. Earlier in the course, we have explored one of them, the WKB approximation, which is adequate for a particle moving through a soft potential profile. In this chapter, we will discuss other techniques that are more suitable for other cases. The historic name for these techniques is the perturbation theory, though it is more fair to speak about several perturbative approaches, because they are substantially different for different cases.

The simplest version of the perturbation theory addresses the problem of stationary states and energy levels of systems described by time-independent Hamiltonians of the type

\[ \hat{H} = \hat{H}^{(0)} + \hat{H}^{(1)}, \]  

where the operator \( \hat{H}^{(1)} \), describing the system’s “perturbation”, is relatively small – in the sense that its addition to the unperturbed operator \( \hat{H}^{(0)} \) results in a relatively small change of the eigenenergies \( E_n \) of the system, and the corresponding eigenstates. A typical problem of this type is the 1D weakly anharmonic oscillator (Fig. 1), described by the Hamiltonian (1) with

\[ \hat{H}^{(0)} = \frac{\hat{p}^2}{2m} + \frac{m \omega_0^2 x^2}{2}, \quad \hat{H}^{(1)} = \alpha x^3 + \beta x^4 + ... \]  

(6.2)

with small coefficients \( \alpha, \beta, \ldots \).

Fig. 6.1. The simplest application of the perturbation theory: a weakly anharmonic 1D oscillator. (Dashed lines characterize the unperturbed, harmonic oscillator.)
I will use this system as our first example, but let me start by describing the perturbative approach to the general time-independent Hamiltonian (1). In the bra-ket formalism, the eigenproblem (4.68) for the perturbed Hamiltonian, i.e. the stationary Schrödinger equation of the system is

\[ (\hat{H}^{(0)} + \hat{H}^{(1)}) |n\rangle = E_n |n\rangle. \]  

(6.3)

Let the eigenstates and eigenvalues of the unperturbed Hamiltonian, which satisfy the equation

\[ \hat{H}^{(0)} |n^{(0)}\rangle = E_n^{(0)} |n^{(0)}\rangle, \]  

(6.4)

be considered as known. In this case, the solution of problem (3) means finding, first, its perturbed eigenvalues \( E_n \) and, second, the coefficients \( \langle n^{(0)} | n' \rangle \) of the expansion of the perturbed state’s vectors \( |n\rangle \) in series over the unperturbed ones, \( |n^{(0)}\rangle \):

\[ |n\rangle = \sum_{n'} |n'^{(0)}\rangle \langle n^{(0)} | n' \rangle. \]  

(6.5)

Let us plug Eq. (5), with the summation index \( n' \) replaced with \( n'' \) (just to have a more compact notation in our forthcoming result), into both sides of Eq. (3):

\[ \sum_{n''} \langle n''^{(0)} | n\rangle \hat{H}^{(0)} |n^{(0)}\rangle + \sum_{n''} \langle n''^{(0)} | n\rangle \hat{H}^{(1)} |n^{(0)}\rangle = \sum_{n''} \langle n''^{(0)} | n\rangle E_n |n^{(0)}\rangle, \]  

(6.6)

and then inner-multiply all terms by an arbitrary unperturbed bra-vector \( \langle n^{(0)} | \) of the system. Assuming that the unperturbed eigenstates are orthonormal, \( \langle n^{(0)} | n''^{(0)} \rangle = \delta_{n'n''} \), and using Eq. (4) in the first term on the left-hand side, we get the following system of linear equations

\[ \sum_{n''} \langle n''^{(0)} | n\rangle H_{nn''}^{(1)} = \langle n^{(0)} | n\rangle (E_n - E_n^{(0)}), \]  

(6.7)

where the matrix elements of the perturbation are calculated, by definition, in the unperturbed brackets:

\[ H_{nn''}^{(1)} = \langle n'^{(0)} | \hat{H}^{(1)} | n''^{(0)} \rangle. \]  

(6.8)

The linear equation system (7) is still exact, and is frequently used for numerical calculations. (Since the matrix coefficients (8) typically decrease when \( n' \) and/or \( n'' \) become sufficiently large, the sum on the left-hand side of Eq. (7) may usually be truncated, still giving an acceptable accuracy of the solution.) To get analytical results, we need to make approximations. In the simple perturbation theory we are discussing now, this is achieved by the expansion of both the eigenenergies and the expansion coefficients into the Taylor series in a certain small parameter \( \mu \) of the problem:

\[ E_n = E_n^{(0)} + E_n^{(1)} + E_n^{(2)} + ..., \]  

(6.9)

\[ \langle n^{(0)} | n\rangle = \langle n^{(0)} | n\rangle^{(0)} + \langle n^{(0)} | n\rangle^{(1)} + \langle n^{(0)} | n\rangle^{(2)} + ..., \]  

(6.10)

where

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1 Please note the similarity of Eq. (7) with Eq. (2.215) of the 1D band theory. Indeed, the latter equation is not much more than a particular form of Eq. (7) for the 1D wave mechanics, and a specific (periodic) potential \( U(x) \) considered as the perturbation Hamiltonian. Moreover, the whole approximate treatment of the weak-potential limit in Sec. 2.7 is essentially a particular case of the perturbation theory we are discussing now (in its 1st order).
\[ E_n^{(k)} \propto \langle n^{(0)} \mid n \rangle^{(k)} \propto \mu^k. \]  

(6.11)

In order to explore the 1\textsuperscript{st}-order approximation, which ignores all terms \(O(\mu^2)\) and higher, let us plug only the two first terms of the expansions (9) and (10) into the basic equation (7):

\[
\sum_{n''} H^{(1)}_{nn''} \left( \delta_{nn''} + \langle n^{(0)} \mid n\rangle^{(1)} \right) = \left( \delta_{nn''} + \langle n^{(0)} \mid n\rangle^{(1)} \right) \left( E_n^{(0)} + E_n^{(1)} - E_{n''}^{(0)} \right). 
\]  

(6.12)

Now let us open the parentheses, and disregard all the remaining terms \(O(\mu^2)\). The result is

\[
H^{(1)}_{nn''} = \delta_{nn''} E_n^{(1)} + \langle n^{(0)} \mid n\rangle^{(1)} (E_n^{(0)} - E_{n''}^{(0)}), \]  

(6.13)

This relation is valid for any set of indices \(n\) and \(n''\); let us start from the case \(n = n''\), immediately getting a very simple (and practically, the most important!) result:

\[
E_n^{(1)} = H^{(1)}_{nn} \equiv \langle n^{(0)} \mid \hat{H}^{(1)} \mid n^{(0)} \rangle. 
\]  

(6.14)

For example, let us see what this result give for two first perturbation terms in the weakly anharmonic oscillator (2):

\[
E_n^{(1)} = \alpha \langle n^{(0)} \mid \hat{x}^3 n^{(0)} \rangle + \beta \langle n^{(0)} \mid \hat{x}^4 n^{(0)} \rangle. 
\]  

(6.15)

As the reader knows (or should know :-) from the solution of Problem 5.9, the first bracket equals zero, while the second one yields\(^2\)

\[
E_n^{(1)} = \frac{3}{4} \beta x_0^4 \left( 2n^2 + 2n + 1 \right). 
\]  

(6.16)

Naturally, there should be some non-vanishing contribution to the energies from the (typically, larger) perturbation proportional to \(\alpha\), so that for its calculation we need to explore the 2\textsuperscript{nd} order of the theory. However, before doing that, let us complete our discussion of its 1\textsuperscript{st} order.

For \(n' \neq n\), Eq. (13) may be used to calculate the eigenstates rather than the eigenvalues:

\[
\langle n''^{(0)} \mid n \rangle^{(1)} = \frac{H^{(1)}_{nn''}}{E_n^{(0)} - E_{n''}^{(0)}}, \quad \text{for } n' \neq n. 
\]  

(6.17)

This means that the eigenket’s expansion (5), in the 1\textsuperscript{st} order, may be represented as

\[
\left| n^{(1)} \right> = C \left| n^{(0)} \right> + \sum_{n' \neq n} \frac{H^{(1)}_{nn'}}{E_n^{(0)} - E_{n'}^{(0)}} \left| n''^{(0)} \right>. 
\]  

(6.18)

The coefficient \(C \equiv \langle n^{(0)} \mid n^{(1)} \rangle\) cannot be found from Eq. (17); however, requiring the final state \(n\) to be normalized, we see that other terms may provide only corrections \(O(\mu^2)\), so that in the 1\textsuperscript{st} order we should take \(C = 1\). The most important feature of Eq. (18) is its denominators: the closer are the unperturbed eigenenergies of two states, the larger is their mutual “interaction” due to the perturbation.

\(^2\) A useful exercise for the reader: analyze the relation between Eq. (16) and the result of the classical theory of such weakly anharmonic (“nonlinear”) oscillator – see, e.g., CM Sec. 5.2, in particular, Eq. (5.49).
This feature also affects the 1st-order’s validity condition, which may be quantified using Eq. (17): the magnitudes of the brackets it describes have to be much less than the unperturbed bracket \( \langle n|n\rangle^{(0)} = 1 \), so that all elements of the perturbation matrix have to be much less than the difference between the corresponding unperturbed energies. For the anharmonic oscillator’s energy corrections (16), this requirement is reduced to \( E_n^{(1)} \ll \hbar \omega_0 \).

Now we are ready for going after the 2nd-order approximation to Eq. (7). Let us focus on the case \( n' = n \), because as we already know, only this term will give us a correction to the eigenenergies. Moreover, since the left-hand side of Eq. (7) already has a small factor \( \hat{H}^{(1)}_{n,n'} \propto \mu \), the bracket coefficients in that part may be taken from the 1st-order result (17). As a result, we get

\[
E_n^{(2)} = \sum_{n'} \langle n|n^{(0)}\rangle \langle n^{(0)}|n\rangle^{(1)} H_{n,n'}^{(1)} = \sum_{n'\neq n} \frac{H_{n,n'}^{(1)}}{E_n^{(0)} - E_{n'}^{(0)}}. \tag{6.19}
\]

Since \( \hat{H}^{(1)} \) has to be Hermitian, we may rewrite this expression as

\[
E_n^{(2)} = \sum_{n'\neq n} \frac{|H_{n,n'}^{(1)}|^2}{E_n^{(0)} - E_{n'}^{(0)}} \equiv \sum_{n'\neq n} \frac{|\langle n|n^{(0)}\rangle \langle n^{(0)}|\hat{H}^{(1)}|n\rangle^{(0)}|^2}{E_n^{(0)} - E_{n'}^{(0)}}. \tag{6.20}
\]

This is the much-celebrated 2nd-order perturbation result, which frequently (in sufficiently symmetric problems) is the first non-vanishing correction to the state energy – for example, from the cubic term (proportional to \( \alpha \)) in our weakly anharmonic oscillator problem (2). To calculate the corresponding correction, we may use another result of the solution of Problem 5.9:

\[
\langle n'\rangle^3 = \left( \frac{x_0}{\sqrt{2}} \right)^3 \times \left\{ n(n-1)(n-2) \right\}^{1/2} \delta_{n,n-3} + 3n^{3/2} \delta_{n,n-1} + 3(n+1)^{1/2} \delta_{n,n+1} + \left[ (n+1)(n+2)(n+3) \right]^{1/2} \delta_{n,n+3}. \tag{6.21}
\]

So, according to Eq. (20), we need to calculate

\[
E_n^{(2)} = \alpha^3 \left( \frac{x_0}{\sqrt{2}} \right)^6 \times \sum_{n'\neq n} \frac{\left\{ n(n-1)(n-2) \right\}^{1/2} \delta_{n,n-3} + 3n^{3/2} \delta_{n,n-1} + 3(n+1)^{1/2} \delta_{n,n+1} + \left[ (n+1)(n+2)(n+3) \right]^{1/2} \delta_{n,n+3} }{\hbar \omega_0 (n-n')} \tag{6.22}
\]

The summation is not as cumbersome as may look, because at the curly brackets’ squaring, all mixed products are proportional to the products of different Kronecker deltas and hence vanish, so that we need to sum up only the squares of each term, finally getting

\[
E_n^{(2)} = -\frac{15}{4} \frac{\alpha^2 x_0^6}{\hbar \omega_0} \left( n^2 + n + \frac{11}{30} \right). \tag{6.23}
\]

This formula shows that all energy level corrections are negative, regardless of the sign of \( \alpha \).\textsuperscript{3} On the contrary, the 1st order correction \( E_n^{(1)} \), given by Eq. (16), does depend on the sign of \( \beta \), so that the net correction, \( E_n^{(1)} + E_n^{(2)} \), may be of any sign.

\textsuperscript{3} Note this is correct for the ground-state energy correction \( E_g^{(2)} \) of any system, because for this state, the denominators of all terms of the sum (20) are negative, while their numerators are always non-negative.
The results (18) and (20) are clearly inapplicable to the degenerate case where, in the absence of perturbation, several states correspond to the same energy level, because of the divergence of their denominators.\(^4\) This divergence hints that in this case, the largest effect of the perturbation is the \textit{degeneracy lifting}, e.g., splitting of the initially degenerate energy level \(E^{(0)}\) (Fig. 2), and that for the analysis of this case we can, to the first approximation, ignore the effect of all other energy levels. (A careful analysis shows that this is indeed the case until the level splitting becomes comparable with the distance to other energy levels.)

Limiting the summation in Eq. (7) to the group of \(N\) degenerate states with equal \(E^{n'}_n(0) \equiv E^{(0)}\), we reduce it to

\[
\sum_{n=1}^{N} \langle n''(0) | n \rangle H^{(1)}_{n'n} - \langle n''(0) | n \rangle (E^{(0)}_n - E^{(0)}),
\]

where now \(n'\) and \(n''\) number the \(N\) states of the degenerate group.\(^5\) For \(n = n'\), Eq. (24) may be rewritten as

\[
\sum_{n=1}^{N} \left( H^{(1)}_{n'n} - E^{(1)}_n \delta_{n'n} \right) n''(0) | n'\rangle = 0, \quad \text{where} \quad E^{(1)}_n \equiv E_n - E^{(0)}.
\]

For each \(n' = 1, 2, \ldots N\), this is a system of \(N\) linear, homogeneous equations (with \(N\) terms each) for \(N\) unknown coefficients \(n''(0) | n'\rangle\). In this problem, we may readily recognize the problem of diagonalization of the perturbation matrix \(H^{(1)}\) – cf. Sec. 4.4 and in particular Eq. (4.101). As in the general case, the condition of self-consistency of the system is:

\[
\begin{vmatrix}
H^{(1)}_{11} - E^{(1)}_n & H^{(1)}_{12} & \cdots \\
H^{(1)}_{21} & H^{(1)}_{22} - E^{(1)}_n & \cdots \\
\vdots & \vdots & \ddots
\end{vmatrix} = 0,
\]

where now the index \(n\) numbers the \(N\) roots of this equation, in an arbitrary order. According to the definition (25) of \(E^{(1)}_n\), the resulting \(N\) energy levels \(E_n\) may be found as \(E^{(0)} + E^{n(t)}_n\). If the perturbation matrix is diagonal in the chosen basis \(n''(0)\), the result is extremely simple,

\[
E_n - E^{(0)} \equiv E^{n(t)}_n = H^{(1)}_{nn},
\]

\(^4\) This is exactly the reason why such simple perturbation approach runs into serious problems for systems with a continuous spectrum, and other techniques (such as the WKB approximation) are often necessary.

\(^5\) Note that here the choice of the basis is to some extent arbitrary, because due to the linearity of equations of quantum mechanics, any linear combination of the states \(n''(0)\) is also an eigenstate of the unperturbed Hamiltonian. However, for using Eq. (25), these combinations have to be orthonormal, as was supposed at the derivation of Eq. (7).
and formally coincides with Eq. (14) for the non-degenerate case, but now it may give a different result for each of \( N \) previously degenerate states \( n \).

Let us see what does this theory give for several important examples. First of all, let us consider a system with two degenerate states with energy sufficiently far from all other levels. Then, in the basis of these two degenerate states, the most general perturbation matrix is

\[
H^{(1)} = \begin{pmatrix} H_{11} & H_{12} \\ H_{21} & H_{22} \end{pmatrix}
\]  

(6.28)

This matrix coincides with the general matrix (5.2) of a two-level system. Hence, we come to the very important conclusion: for a weak perturbation, all properties of any double-degenerate system are identical to those of the genuine two-level systems, which were the subject of numerous discussions in Chapter 4 and again in Sec. 5.1. In particular, its eigenenergies are given by Eq. (5.6), and may be described by the level-anticrossing diagram shown in Fig. 5.1.

**6.2. The linear Stark effect**

As a more involved example of the level degeneracy lifting by a perturbation, let us discuss the Stark effect\(^6\) – the atomic level splitting by an external electric field. Let us study this effect, in the linear approximation, for a hydrogen-like atom/ion. Taking the direction of the external electric field \( \vec{E} \) (which is practically always uniform on the atomic scale) for the \( z \)-axis, the perturbation may be represented by the following Hamiltonian:

\[
\hat{H}^{(1)} = -F\hat{\vec{z}} = -q\vec{E} \hat{\vec{z}} = -qE \cos \theta.
\]  

(6.29)

(In the last form, the operator sign is dropped, because we will work in the coordinate representation.)

As you (should :-) remember, energy levels of a hydrogen-like atom depend only on the principal quantum number \( n \) – see Eqs. (1.12) or (3.201); hence all the states, besides the ground state \( n = 1 \) (“1s” in the spectroscopic nomenclature) in which \( l = m = 0 \), have some degeneracy, which grows rapidly with \( n \). I will carry out the calculations only for the lowest degenerate level with \( n = 2 \). Since, according to Eq. (3.203), \( 0 \leq l \leq n - 1 \), at this level the orbital quantum number \( l \) may equal either 0 (one \( 2s \) state, with \( m = 0 \)) or 1 (three \( 2p \) states, with \( m = 0, \pm 1 \)). Due to this 4-fold degeneracy, \( H^{(1)} \) is a \( 4 \times 4 \) matrix with 16 elements:

\[
H^{(1)} = \begin{pmatrix} H_{11} & H_{12} & H_{13} & H_{14} \\ H_{21} & H_{22} & H_{23} & H_{24} \\ H_{31} & H_{32} & H_{33} & H_{34} \\ H_{41} & H_{42} & H_{43} & H_{44} \end{pmatrix}
\]

(6.30)

\( \begin{array}{c|c} l=0 \\ \hline m=0 & m=0 & m=+1 & m=-1 \\ \end{array} \)

\( \begin{array}{c|c} l=1 \\ \hline m=0 & m=0 & m=+1 & m=-1 \\ \end{array} \)

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\(^6\) This effect was discovered experimentally in 1913 by Johannes Stark and (independently) by Antonio Lo Surdo, so it is sometimes (and more fairly) called the “Stark – Lo Surdo effect”. Sometimes this name is used with the qualifier “dc” to distinguish it from the *ac* Stark effect – the energy level shift under the effect of an ac field – see Sec. 5.
However, there is no need to be scared. First, due to the Hermitian nature of the operator, only ten of the matrix elements (four diagonal ones and six off-diagonal elements) may be substantially different. Moreover, due to a high symmetry of the problem, there are a lot of zeros even among these elements. Indeed, let us have a look at the angular components $Y_l^m$ of the corresponding wavefunctions, with $l = 0$ and $l = 1$, described by Eqs. (3.174)-(3.175). For the states with $m = \pm 1$, the azimuthal parts of wavefunctions are proportional to $\exp\{\pm i\varphi\}$; hence the off-diagonal elements $H_{34}$ and $H_{43}$ of the matrix (30), relating these functions, are proportional to

$$\int_{0}^{2\pi} d\varphi \left( e^{\pm i\varphi} \right)^* = 0.$$

(6.31)

The azimuthal-angle symmetry also kills the off-diagonal elements $H_{13}, H_{14}, H_{23}, H_{24}$ (and hence their complex conjugates $H_{31}, H_{41}, H_{32},$ and $H_{42}$), because they relate the states with $m = 0$ and $m \neq 0$, and hence are proportional to

$$\int_{0}^{2\pi} d\varphi e^{\pm i\varphi} = 0.$$

(6.32)

For the diagonal elements $H_{33}$ and $H_{44}$, corresponding to $m = \pm 1$, the azimuthal-angle integral does not vanish, but since the corresponding spherical functions depend on the polar angle as $\sin \theta$, the matrix elements are proportional to

$$\int_{0}^{\pi} \sin \theta d\theta \sin \theta \cos \theta \sin \theta = \int_{-1}^{1} \cos \theta (1 - \cos^2 \theta) d(\cos \theta),$$

(6.33)

and hence are equal to zero – as any limit-symmetric integral of an odd function. Finally, for the states $2s$ and $2p$ with $m = 0$, the diagonal elements $H_{11}$ and $H_{22}$ are also killed by the polar-angle integration:

$$\int_{0}^{\pi} \sin \theta d\theta \cos \theta = \int_{-1}^{1} \cos \theta d(\cos \theta) = 0,$$

(6.34)

and

$$\int_{0}^{\pi} \sin \theta d\theta \cos^3 \theta = \int_{-1}^{1} \cos^3 \theta d(\cos \theta) = 0.$$

(6.35)

Hence, the only non-vanishing elements of the matrix (30) are two off-diagonal elements $H_{12}$ and $H_{21}$, which relate two states with the same $m = 0$, but different $l = 0, 1$, because they are proportional to

$$\int_{0}^{\pi} \sin \theta d\theta \cos^2 \theta = \frac{1}{\sqrt{3}} \neq 0.$$

(6.36)

What remains is to use Eqs. (3.209) for the radial parts of these functions to complete the calculation of those two matrix elements:

$$H_{12} = H_{21} = -\frac{q \hbar}{\sqrt{3}} \int_{0}^{\infty} r^2 dr \mathcal{R}_{2,0}(r)r^2 \mathcal{R}_{2,1}(r).$$

(6.37)

Due to the additive structure of the function $\mathcal{R}_{2,0}(r)$, the integral falls into a sum of two table integrals, both of the type MA Eq. (6.7d), finally giving
\[ H_{12} = H_{21} = 3q \varepsilon r_0, \quad (6.38) \]

where \( r_0 \) is the spatial scale (3.192); for the hydrogen atom, it is just the Bohr radius \( r_B \) (1.10).

Thus, the perturbation matrix (30) is reduced to

\[
H^{(1)} = \begin{pmatrix}
0 & 3q \varepsilon r_0 & 0 & 0 \\
3q \varepsilon r_0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0
\end{pmatrix},
\quad (6.39)
\]

so that the condition (26) of the self-consistency of the system (25),

\[
\begin{pmatrix}
-E_n^{(1)} & 3q \varepsilon r_0 & 0 & 0 \\
3q \varepsilon r_0 & -E_n^{(1)} & 0 & 0 \\
0 & 0 & -E_n^{(1)} & 0 \\
0 & 0 & 0 & -E_n^{(1)}
\end{pmatrix} = 0,
\quad (6.40)
\]

gives a very simple characteristic equation

\[
(E_n^{(1)})^2 - (3q \varepsilon r_0)^2 = 0.
\quad (6.41)
\]

with the roots

\[
E_{1,2}^{(1)} = 0, \quad E_{3,4}^{(1)} = \pm 3q \varepsilon r_0.
\quad (6.42)
\]

so that the degeneracy is only partly lifted – see the levels in Fig. 3.7

\[
|+\rangle = \frac{1}{\sqrt{2}} (|2s\rangle + |2p\rangle) \quad m = 0 \\
E_2^{(0)} \quad 3q \varepsilon r_0 \quad m = \pm 1 \\
|\rangle = \frac{1}{\sqrt{2}} (|2s\rangle - |2p\rangle) \quad m = 0 \quad \text{Fig. 6.3. The linear Stark effect for the level } n = 2 \text{ of a hydrogen-like atom.}
\]

Generally, in order to understand the nature of states corresponding to these levels, we should return to Eq. (25) with each calculated value of \( E_n^{(1)} \), and find the corresponding expansion coefficients \( \langle n'|n\rangle \), which describe the perturbed states. However, in our simple case, the outcome of this procedure is clear in advance. Indeed, since the states with \( |l = 1, m = \pm 1\rangle \) are not affected by the perturbation at all (in the linear approximation in the electric field), their degeneracy is not lifted, and energy is not affected – see the middle line in Fig. 3. On the other hand, the partial perturbation matrix connecting the states \( 2s \) and \( 2p \), i.e. the top left \( 2\times2 \) part of the full matrix (39), is proportional to the

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7 The proportionality of this splitting to the small field is responsible for the qualifier “linear” in the name of this effect. If observable effects grow only as \( \varepsilon^2 \) (see, e.g., Problem 9), the term quadratic Stark effect is used instead.
Pauli matrix $\sigma$, and we already know the result of its diagonalization – see Eqs. (4.113)-(4.114). This means that the upper and lower split levels correspond to very simple linear combinations of the previously degenerate states with $m = 0$,

$$|\pm\rangle = \frac{1}{\sqrt{2}}(|2s\rangle \pm |2p\rangle).$$  \hfill (6.43)

Finally, let us estimate the magnitude of the linear Stark effect for a hydrogen atom. For a very high electric field of $E = 3 \times 10^6$ V/m,$^8$ $|q| = e \approx 1.6 \times 10^{-19}$ C, and $r_0 = r_B \approx 0.5 \times 10^{-10}$ m, we get a level splitting of $3qER_0 \approx 0.8 \times 10^{-22}$ J $\approx 0.5$ meV. This number is much lower than the unperturbed energy of the level, $E_2 = -E_H/(2 \times 2^2) \approx -3.4$ eV, so that the perturbative result is quite applicable. On the other hand, the calculated splitting is much larger than the resolution limit imposed by the natural linewidth ($\sim 10^{-7} E_2$, see Chapter 9), so that the effect is quite observable even in substantially lower electric fields. Note, however, that our simple results are quantitatively correct only when the Stark splitting (42) is much larger than the fine-structure splitting of the same level in the absence of the field-- see the next section.

6.3. Fine structure of atomic levels

Now let us use the same perturbation theory analyze, also for the simplest case of a hydrogen-like atom/ion, the so-called fine structure of atomic levels – their degeneracy lifting even in the absence of external fields. In the limit when the effective speed $v$ of the electron motion is much smaller than the speed of light $c$ (as it is in the hydrogen atom), the fine structure may be analyzed as a sum of two small relativistic effects. To analyze the first of them, let us expand the well-known classical relativistic expression$^9$ for the kinetic energy $T = E - mc^2$ of a free particle with the rest mass $m$,\(^10\)

$$T = \left(\frac{m^2c^4 + p^2c^2}{m^2c^2}\right)^{1/2} - mc^2 \equiv mc^2 \left[1 + \frac{p^2}{m^2c^2}\right]^{1/2} - 1,$$  \hfill (6.44)

into the Taylor series with respect to the small ratio $(p/mc)^2 \approx (v/c)^2$:

$$T = mc^2 \left[1 + \frac{1}{2}\left(\frac{p}{mc}\right)^2 + \frac{1}{8}\left(\frac{p}{mc}\right)^4 + ... - 1\right] \equiv \frac{p^2}{2m} - \frac{p^4}{8m^3c^2} + ...,$$  \hfill (6.45)

and drop all the terms besides the first (non-relativistic) one and the next spelled-out term, which represents the first relativistic correction to $T$.

Following the correspondence principle, the quantum-mechanical problem in this approximation may be described by the perturbative Hamiltonian (1), where the unperturbed (non-relativistic) Hamiltonian of the problem (whose eigenstates and eigenenergies were discussed in Sec. 3.5) is

$$\hat{H}^{(0)} = \frac{\hat{p}^2}{2m} + \hat{U}(r), \quad \hat{U}(r) = -\frac{C}{r},$$  \hfill (6.46)

$^8$ This value approximately corresponds to the threshold of electric breakdown in air at ambient conditions, due to the impact ionization. As a result, experiments with higher dc fields are rather difficult.

$^9$ See, e.g., EM Eq. (9.78) – or any undergraduate text on the special relativity.

$^{10}$ This fancy font is used, as in Secs. 3.5-3.8, to distinguish the mass $m$ from the magnetic quantum number $m$. 

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while the kinetic-relativistic perturbation is

$$\hat{H}^{(1)} = -\frac{\hat{p}^4}{8m^3c^2} \equiv -\frac{1}{2mc^2}\left(\frac{\hat{p}^2}{2m}\right)^2. \quad (6.47)$$

Using Eq. (46), we may rewrite the last formula as

$$\hat{H}^{(1)} = -\frac{1}{2mc^2}\left(\hat{H}^{(0)} - \hat{U}(r)\right)^2, \quad (6.48)$$

so that its matrix elements participating in the characteristic equation (25) for a given degenerate energy level (3.201), i.e. a given principal quantum number $n$, are

$$\langle nlm | \hat{H}^{(1)} | n'l'm' \rangle = -\frac{1}{2mc^2} \langle nlm | \left(\hat{H}^{(0)} - \hat{U}(r)\right)^2 | n'l'm' \rangle, \quad (6.49)$$

where the bra- and ket-vectors describe the unperturbed eigenstates, whose eigenfunctions (in the coordinate representation) are given by Eq. (3.200): $\psi_{n,l,m} = \mathcal{R}_n(r)Y_{l,m}^m(\theta, \phi)$.

It is straightforward (and hence left for the reader ;-) to prove that all off-diagonal elements of the set (49) are equal to 0. Thus we may use Eq. (27) for each set of the quantum numbers $\{n, l, m\}$:

$$E^{(1)}_{n,l,m} \equiv E^{(0)}_{n,l,m} - E^{(0)}_n = \langle nlm | \hat{H}^{(1)} | nlm \rangle = -\frac{1}{2mc^2} \left(\hat{H}^{(0)} - \hat{U}(r)\right)^2|_{n,l,m}$$

$$= -\frac{1}{2mc^2} \left(E_n^2 - 2E_n\langle \hat{U} \rangle_{n,l} + \langle \hat{U}^2 \rangle_{n,l} \right) = -\frac{1}{2mc^2} \left(\frac{E_0^2}{4n^2} - \frac{E_0}{n^2} C \left(\frac{1}{r}\right)^{n,l} + C^2 \left(\frac{1}{r^2}\right)^{n,l} \right), \quad (6.50)$$

where the index $m$ has been dropped, because the radial wavefunctions $\mathcal{R}_n(r)$, which affect these expectation values, do not depend on that quantum number. Now using Eqs. (3.191), (3.201) and the first two of Eqs. (3.211), we finally get

$$E^{(1)}_{n,l} = -\frac{mC^2}{2\hbar^2c^2n^4} \left(\frac{n}{l+\frac{1}{2}} - \frac{3}{4}\right) \equiv -\frac{2E_n^2}{mc^2} \left(\frac{n}{l+\frac{1}{2}} - \frac{3}{4}\right). \quad (6.51)$$

Let us discuss this result. First of all, its last form confirms that the correction (51) is indeed much smaller than the unperturbed energy $E_n$ (and hence the perturbation theory is valid) if the latter is much smaller than the relativistic rest energy $mc^2$ of the particle – as it is for the hydrogen atom. Next, since in the Bohr problem’s solution $n \geq l + 1$, the first fraction in the parentheses of Eq. (51) is always larger than 1, and hence than 3/4, so that the kinetic relativistic correction to energy is negative for all $n$ and $l$. (Actually, this fact could be predicted already from Eq. (47), which shows that the perturbation’s Hamiltonian is a negatively defined form.) Finally, for a fixed principal number $n$, the negative correction’s magnitude decreases with the growth of $l$. This fact may be interpreted using the second of Eqs. (3.211): the larger is $l$ (at fixed $n$), the larger is the particle’s effective distance from the center, and hence the smaller is its effective velocity, i.e. the smaller is the magnitude of the quantum-mechanical average of the negative relativistic correction (47) to the kinetic energy.

The result (51) is valid for the Coulomb interaction $U(r) = -C/r$ of any physical nature. However, if we speak specifically about hydrogen-like atoms/ions, there is also another relativistic correction to energy, due to the so-called spin-orbit interaction (alternatively called the “spin-orbit coupling”). Its physics may be understood from the following semi-quantitative classical reasoning: from the “the point...
of view" of an electron rotating about the nucleus at constant distance $r$ with velocity $v$, it is the nucleus, of the electric charge $Ze$, that rotates about the electron with the velocity $(-v)$ and hence the time period $T = 2\pi r/v$. From the point of view of magnetostatics, such circular motion of the electric charge $Q = Ze$, is equivalent to a circular dc electric current $I = Q/T = (Ze)(v/2\pi)$, which creates, at the electron’s location, i.e. in the center of the current loop, the magnetic field with the following magnitude:$^{11}$

$$B = \frac{\mu_0}{2r} I \equiv \frac{\mu_0 Ze}{2r 2\pi} \equiv \frac{\mu_0 Ze}{4\pi r^2}.$$  (6.52)

The field’s direction $\mathbf{n}$ is perpendicular to the apparent plane of the nucleus’ rotation (i.e. that of the real rotation of the electron), and hence its vector may be readily expressed via the similarly directed vector $\mathbf{L} = m_e vr \mathbf{n}$ of the electron’s angular (orbital) momentum:

$$B = \frac{\mu_0}{4\pi^2} \frac{Ze}{m_e} m_e v r \mathbf{n} \equiv \frac{\mu_0}{4\pi^2} \frac{Ze}{m_e} \mathbf{L} \equiv \frac{Ze}{4\pi \varepsilon_0 r^3 m_e c^2} \mathbf{L},$$  (6.53)

where the last step used the basic relation between the SI-unit constants: $\mu_0 \equiv 1/c^2 \varepsilon_0$.

A more careful (but still classical) analysis of the problem$^{12}$ brings both good and bad news. The bad news is that the result (53) is wrong by the so-called Thomas factor of two even for the circular motion, because the electron moves with acceleration, and the reference frame bound to it cannot be considered inertial (as was implied in the above reasoning), so that the effective magnetic field felt by the electron is actually

$$\mathcal{B} = \frac{Ze}{8\pi \varepsilon_0 r^3 m_e c^2} \mathbf{L}.$$  (6.54)

The good news is that, so corrected, the result is valid not only for circular but for an arbitrary orbital motion in the Coulomb field $U(r)$. Hence from the discussion in Sec. 4.1 and Sec. 4.4 we may expect that the quantum-mechanical description of the interaction between this effective magnetic field and the electron’s spin moment (4.115) is given by the following perturbation Hamiltonian$^{13}$

$$\hat{H}^{(1)} = -\hat{\mathbf{m}} \cdot \mathcal{B} = -\gamma_e \hat{\mathbf{S}} \left( \frac{Ze}{8\pi \varepsilon_0 r^3 m_e c^2} \hat{\mathbf{L}} \right) \equiv \frac{1}{2m_e c^2} \frac{Ze^2}{4\pi \varepsilon_0 r^3} \frac{1}{2} \hat{\mathbf{S}} \cdot \hat{\mathbf{L}},$$  (6.55)

where at spelling out the electron’s gyromagnetic ratio $\gamma_e = -g_e e/2m_e$, the small correction to the value $g_e = 2$ of the electron’s $g$-factor (see Sec. 4.4) has been ignored, because Eq. (55) is already a small correction. This expression is confirmed by the fully-relativistic Dirac theory, to be discussed in Sec. 9.7 below: it yields, for an arbitrary central potential $U(r)$, the following spin-orbit coupling Hamiltonian:

$$\hat{H}^{(1)} = \frac{1}{2m_e c^2} \frac{1}{r} \frac{dU(r)}{dr} \hat{\mathbf{S}} \cdot \hat{\mathbf{L}}.$$  (6.56)

$^{11}$ See, e.g., EM Sec. 5.1, in particular, Eq. (5.24). Note that such effective magnetic field is induced by any motion of electrons, in particular that in solids, leading to a variety of spin-orbit level-splitting effects there – see, e.g., a concise review by R. Winkler et al., in B. Kramer (ed.), Advances in Solid State Physics 41, 211 (2001).

$^{12}$ It was carried out first by Llewellyn Thomas in 1926; for a simple review see, e.g., R. Harr and L. Curtis, Am. J. Phys. 55, 1044 (1987).

$^{13}$ In the Gaussian units, Eq. (55) is valid without the factor $4\pi \varepsilon_0$ in the denominator; while Eq. (56), “as is”.  

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Spin-orbit coupling
For the Coulomb potential $U(r) = -\frac{Ze^2}{4\pi\varepsilon_0 r}$, this formula is reduced to Eq. (55).

As we already know from the discussion in Sec. 5.7, the angular factor of this Hamiltonian commutes with all the operators of the coupled-representation group (inside the blue line in Fig. 5.12): $\hat{L}^2$, $\hat{S}^2$, $\hat{J}^2$, and $\hat{J}_z$, and hence is diagonal in the coupled-representation basis with definite quantum numbers $l$, $j$, and $m_j$ (and of course $s = \frac{1}{2}$). Hence, using Eq. (5.181) to rewrite Eq. (56) as

$$\hat{H}^{(i)} = \frac{1}{2m_e^2c^2} \frac{Ze^2}{4\pi\varepsilon_0} \frac{1}{r^3} \left( \hat{j}^2 - \hat{L}^2 - \hat{S}^2 \right),$$

(6.57)

we may again use Eq. (27) for each set $\{s, l, j, m_j\}$, with the common $n$:

$$E_{n,j,l}^{(i)} = \frac{1}{2m_e^2c^2} \frac{Ze^2}{4\pi\varepsilon_0} \left( \frac{1}{r^3} \right)_{n,l} \frac{1}{2} \left\langle \hat{j}^2 - \hat{L}^2 - \hat{S}^2 \right\rangle_{j,s},$$

(6.58)

where the indices irrelevant for each particular factor have been dropped. Now using the last of Eqs. (3.211), and similar expressions (5.169), (5.175), and (5.177) for eigenvalues of the involved operators, we get an explicit expression for the spin-orbit corrections $^{14}$

$$E_{n,j,l}^{(i)} = \frac{1}{2m_e^2c^2} \frac{Ze^2}{4\pi\varepsilon_0} \frac{\hbar^2}{2r_0^3} \frac{j(j+1)-l(l+1)-\frac{3}{4}}{n^3l(l+1/2)(l+1)} \equiv \frac{E_z^2}{m_e^2c^2} \frac{n(j(j+1)-l(l+1)-\frac{3}{4})}{l(l+1/2)(l+1)},$$

(6.59)

with $l$ and $j$ related by Eq. (5.189): $j = l \pm \frac{1}{2}$.

The last form of its result shows clearly that this correction has the same scale as the kinetic correction (51). In the 1st order of the perturbation theory, they may be just added (with $m = m_e$), giving a surprisingly simple formula for the net fine structure of the $n$th energy level:

$$E_{\text{fine}}^{(i)} = \frac{E_n^2}{2m_e^2c^2} \frac{3}{j+1/2} \left( 3 - \frac{4n}{j+1/2} \right).$$

(6.60)

This simplicity, as well as the independence of the result of the orbital quantum number $l$, will become less surprising when (in Sec. 9.7) we see that this formula follows in one shot from the Dirac theory, in which the Bohr atom’s energy spectrum in numbered only with $n$ and $j$, but not $l$. Let us recall that for an electron ($s = \frac{1}{2}$), according to Eq. (5.189) with $0 \leq l \leq n - 1$, the quantum number $j$ may take $n$ positive half-integer values, from $\frac{1}{2}$ to $n - \frac{1}{2}$. Hence, Eq. (60) shows that the fine structure of the $n$th Bohr’s energy level has $n$ sub-levels – see Fig. 4.

Please note that according to Eq. (5.175), each of these sub-levels is still $(2j+1)$-times degenerate in the quantum number $m_j$. This degeneracy is very natural, because in the absence of an external field the system is still isotropic. Moreover, on each fine-structure level, besides the highest one ($j = n - \frac{1}{2}$), each of the $m_j$-states is doubly-degenerate in the orbital quantum number $l = j \mp \frac{1}{2}$ – see the

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$^{14}$ The factor $l$ in the denominator does not give a divergence at $l = 0$, because in this case $j = s = \frac{1}{2}$, so that $j(j+1) = \frac{3}{4}$, and the numerator turns into 0 as well. A careful analysis of this case (which may be found, e.g., in G. Woolgate, *Elementary Atomic Structure*, 2nd ed., Oxford, 1983), as well as the exact analysis of the hydrogen atom using the Dirac theory (see Sec. 9.7), show that Eq. (60), which does not include $l$, is valid even in this case.

$^{15}$ This is natural, because the magnetic interaction of charged particles is essentially a relativistic effect, of the same order ($-\gamma^2/c^2$) as the kinetic correction (47) – see, e.g., EM Sec. 5.1, in particular Eq. (5.5).
labels of \( l \) in Fig. 4. (According to Eq. (5.190), each of these states, with fixed \( j \) and \( m_j \), may be represented as a linear combination of two states with adjacent values of \( l \), and hence different electron spin orientations, \( m_s = \pm \frac{1}{2} \), weighed with the Clebsch-Gordan coefficients.)

These details aside, one may crudely say that the relativistic corrections combined make the total eigenenergy grow with \( l \), contributing to the effect already mentioned at our analysis of the periodic table of elements in Sec. 3.7. The relative scale of this increase may be scaled by the largest deviation from the unperturbed energy \( E_n \), reached for \( s \)-states (with \( l = 0, j = \frac{1}{2} \)):

\[
\frac{E_{\text{max}}^{(1)}}{E_n} = \frac{E_n}{m_e c^2} \left(2n - \frac{3}{2}\right) \equiv \left(\frac{Ze^2}{4\pi\varepsilon_0 \hbar c}\right)^2 \left(\frac{1}{n} - \frac{3}{4n^2}\right) \equiv Z^2 \alpha^2 \left(\frac{1}{n} - \frac{3}{4n^2}\right). \tag{6.61}
\]

where \( \alpha \) is the fine-structure (“Sommerfeld’s”) constant,

\[
\alpha \equiv \frac{e^2}{4\pi\varepsilon_0 \hbar c} \approx \frac{1}{137}, \tag{6.62}
\]

(already mentioned in Sec. 4.4), which characterizes the relative strength (or rather weakness :-)) of the electromagnetic effects in quantum mechanics – which in particular makes the perturbative quantum electrodynamics possible.\(^{16}\) These expressions show that the fine-structure splitting is a very small effect (~\( \alpha^2 \sim 10^{-6} \)) for the hydrogen atom, but it rapidly grows (as \( Z^2 \)) with the nuclear charge (i.e. the atomic number) \( Z \), and becomes rather substantial for the heaviest stable atoms with \( Z \sim 10^2 \).

### 6.4. The Zeeman effect

Now, we are ready to review the Zeeman effect – the atomic level splitting by an external magnetic field.\(^{17}\) Using Eq. (3.26), with \( q = -e \), for the description of the electron’s orbital motion in the field, and the Pauli Hamiltonian (4.163), with \( \gamma = -e/m_e \), for the electron spin’s interaction with the field, we see that even for a hydrogen-like (i.e. single-electron) atom/ion, neglecting the relativistic effects, the full Hamiltonian is rather involved:

\[
\hat{H} = \frac{1}{2m_e} \left(\hat{p} + e\hat{A}\right)^2 - \frac{Ze^2}{4\pi\varepsilon_0 r} + \frac{e}{m_e} \mathcal{B} \cdot \hat{S}. \tag{6.63}
\]

\(^{16}\) The expression \( \alpha^2 = E_{\text{H}}/m_e c^2 \), where \( E_{\text{H}} \) is the Hartree energy (1.13), i.e. the scale of energies \( E_n \), is also very revealing.

\(^{17}\) It was discovered experimentally in 1896 by Pieter Zeeman who, amazingly, was fired from the University of Leiden for unauthorized use of lab equipment for this work – just to receive a Nobel Prize for it in a few years!
There are several simplifications we may make. First, let us assume that the external field is spatial-uniform on the atomic scale (which is a very good approximation for most cases), so that we can take its vector potential in an axially-symmetric gauge – cf. Eq. (3.132):

$$A = \frac{1}{2} \mathbf{B} \times \mathbf{r}.$$  (6.64)

Second, let us neglect the terms proportional to $\mathbf{B}^2$, which are small in practical magnetic fields of the order of a few teslas. The remaining term in the effective kinetic energy, describing the interaction with the magnetic field, is linear in the momentum operator, so that we may repeat the standard classical calculation to reduce it to the product of $\mathbf{B}$ by the orbital magnetic moment’s component $m_z = -eL_z/2m_e$ – besides that both $m_z$ and $L_z$ should be understood as operators now. As a result, the Hamiltonian (63) reduces to Eq. (1), $\hat{H}^{(0)} + \hat{H}^{(1)}$, where $\hat{H}^{(0)}$ is that of the atom at $\mathbf{B} = 0$, and

$$\hat{H}^{(1)} = \frac{e\mathbf{B}}{2m_e}(\hat{L}_z + 2\hat{S}_z).$$  (6.65)

This expression immediately reveals the major complication with the Zeeman effect’s analysis. Namely, in comparison with the equal orbital and spin contributions to the total angular momentum (5.171) of the electron, the spin produces a twice larger contribution to the magnetic moment, so that the right-hand side of Eq. (65) is not proportional to the total angular moment $\mathbf{J}$. As a result, the effect’s description is simple only in two limits.

If the magnetic field is so high that its effects are much stronger than the relativistic (fine-structure) effects discussed in the previous section, we may treat the two terms in Eq. (65) as independent perturbations of different (orbital and spin) degrees of freedom. Since each of the perturbation matrices is diagonal in its own $z$-basis, we can again use Eq. (27) to write

$$E - E^{(0)} = \frac{e\mathbf{B}}{2m_e}\left(\langle n,l,m_l|\hat{L}_z|n,l,m_l\rangle + 2\langle m_z|\hat{S}_z|m_z\rangle\right) = \frac{e\mathbf{B}}{2m_e}(\hbar m_l + 2\hbar m_s) = \mu_B \mathbf{B}(m_l \pm 1).$$  (6.66)

This result describes the splitting of each $2(2l + 1)$-degenerate energy level, with certain $n$ and $l$, into $(2l + 3)$ levels (Fig. 5), with the adjacent level distance of $\mu_B \mathbf{B}$, of the order of $\sim 10^{-23}$ J $\sim 10^{-4}$ eV per tesla. Note that all the levels, besides the top and bottom ones, remain doubly degenerate. This limit of the Zeeman effect is sometimes called the Paschen-Back effect – whose simplicity was recognized only in the 1920s, due to the need in very high magnetic fields for its observation.

In the opposite limit of relatively low magnetic fields, the Zeeman effect takes place on the background of the much larger fine-structure splitting. As was discussed in Sec. 3, at $\mathbf{B} = 0$ each split sub-level has a $2(2l + 1)$-fold degeneracy corresponding to $(2j + 1)$ different values of the half-integer quantum number $m_j$, ranging from $-j$ to $+j$, and two values of the integer $l = j \mp \frac{1}{2}$ – see Fig. 4. The

18 Despite its smallness, the quadratic term is necessary for a description of the negative contribution of the orbital motion to the magnetic susceptibility $\chi_m$ (the so-called orbital diamagnetism, see EM Sec. 5.5), whose analysis, using Eq. (63), is left for the reader’s exercise.

19 See, e.g., EM Sec. 5.4, in particular Eqs. (5.95) and (5.100).

20 In the almost-hydrogen-like, but more complex atoms (such as those of alkali metals), the degeneracy in $l$ may be lifted by the electron-electron Coulomb interaction even in the absence of the external magnetic field.
magnetic field lifts this degeneracy. Indeed, in the coupled representation discussed in Sec. 5.7, the perturbation (65) is described by the matrix with elements

\[ H^{(1)} = \frac{eB}{2m_e} \left( j, m_j \right | \hat{L}_z + 2 \hat{S}_z | j', m_j' \right) \equiv \frac{eB}{2m_e} \left( j, m_j \right | \hat{J}_z + \hat{S}_z | j', m_j' \right) \]

\[ = \frac{eB}{2m_e} \left( \hbar m_j \delta_{m_j m_j'} + \left( j, m_j \right | \hat{S}_z | j', m_j' \right) \right) \quad (6.67) \]

To spell out the second term, let us use the general expansion (5.183) for the particular case \( s = \frac{1}{2} \), when (as was discussed in the end of Sec. 5.7) it has at most two non-vanishing terms, with the Clebsh-Gordan coefficients (5.190):

\[
\begin{align*}
| j = l \pm \frac{1}{2}, m_j \rangle &= \pm \left( \frac{l \pm m_j + \frac{1}{2}}{2l + 1} \right)^{1/2} \left| m_j = m_j - \frac{1}{2}, m_s = \mp \frac{1}{2} \right> + \left( \frac{l \mp m_j + \frac{1}{2}}{2l + 1} \right)^{1/2} \left| m_j = m_j + \frac{1}{2}, m_s = \pm \frac{1}{2} \right>. 
\end{align*}
\]

(6.68)

Taking into account that the operator \( \hat{S}_z \) gives non-zero brackets only for \( m_s = m_s' \), the 2x2 matrix of elements \( \left< m_i = m_j \pm \frac{1}{2}, m_s = \mp \frac{1}{2} \right| \hat{S}_z \left| m_i = m_j \pm \frac{1}{2}, m_s = \mp \frac{1}{2} \right> \) is diagonal, so we may use Eq. (27) to get

\[
E - E^{(0)} = \frac{eB}{2m_e} \left[ \hbar m_j + \hbar \left( \frac{l \pm m_j + 1/2}{2l + 1} \right) - \hbar \left( \frac{l \mp m_j + 1/2}{2l + 1} \right) \right] 
\]

\[
\equiv \frac{eB}{2m_e} \hbar m_j \left( 1 \pm \frac{1}{2l + 1} \right) = \mu_B B m_j \left( 1 \pm \frac{1}{2l + 1} \right), \quad \text{for } -j \leq m_j \leq +j, 
\]

(6.69)

where the two signs correspond to the two possible values of \( l = j \pm \frac{1}{2} \) – see Fig. 6.

Anomalous Zeeman effect for \( s = \frac{1}{2} \)

Fig. 6.5. The Paschen-Back effect.

Fig. 6.6. The anomalous Zeeman effect in a hydrogen-like atom/ion.
We see that the magnetic field splits each sub-level of the fine structure, with a given \( l \), into \( 2j + 1 \) equidistant levels, with the distance between the levels depending on \( l \). In the late 1890s, when the Zeeman effect was first observed, there was no notion of spin at all, so that this puzzling result was called the anomalous Zeeman effect. (In this terminology, the normal Zeeman effect is the one with no spin splitting, i.e. without the second terms in the parentheses of Eqs. (66), (67), and (69); it was first observed in 1898 by Preston Thomas in atoms with zero net spin.)

The strict quantum-mechanical analysis of the anomalous Zeeman effect for arbitrary \( s \) (which is important for applications to multi-electron atoms) is conceptually not complex, but requires explicit expressions for the corresponding Clebsch-Gordan coefficients, which are rather bulky. Let me just cite the unexpectedly simple result of this analysis:

\[
\Delta E = \mu_B B m_j g,
\]

(6.70a)

where \( g \) is the so-called Lande factor:

\[
g = 1 + \frac{j(j+1) + s(s+1) - l(l+1)}{2j(j+1)}.
\]

(6.70b)

For \( s = \frac{1}{2} \) (and hence \( j = l \pm \frac{1}{2} \)), this factor is reduced to the parentheses in the last forms of Eq. (69).

It is remarkable that Eqs. (70) may be readily derived using very plausible classical arguments, similar to those used in Sec. 5.7 – see Fig. 5.13 and its discussion. As was discussed in Sec. 5.6, in the absence of spin, the quantization of the observable \( L_z \) is an extension of the classical picture of the torque-induced precession of the vector \( \mathbf{L} \) about the magnetic field’s direction, so that the interaction energy, proportional to \( \mathbf{B} \cdot \mathbf{L} \), remains constant – see Fig. 7a. On the other hand, at the spin-orbit interaction without an external magnetic field, the Hamiltonian function of the system includes the product \( \mathbf{S} \cdot \mathbf{L} \), so that in the stationary state it has to be constant, together with \( J^2 \), \( L^2 \), and \( S^2 \). Hence, this system’s classical image is a precession of the vectors \( \mathbf{S} \) and \( \mathbf{L} \) about the direction of the vector \( \mathbf{J} = \mathbf{L} + \mathbf{S} \), in a manner that the spin-orbit interaction energy, proportional to the product \( \mathbf{L} \cdot \mathbf{S} \), remains constant (Fig. 7b). On this backdrop, the anomalous Zeeman effect in a relatively weak magnetic field \( \mathbf{B} = B \mathbf{n}_z \) corresponds to a much slower precession of the vector \( \mathbf{J} \) (“dragging” the rapidly rotating vectors \( \mathbf{L} \) and \( \mathbf{S} \) with it) about the \( z \)-axis.

**Fig. 6.7.** Classical images of (a) the orbital angular momentum’s quantization in a magnetic field, and (b) the fine-structure level splitting.

\(^{21}\) This formula is frequently used with capital letters \( J, S, \) and \( L \), which denote the quantum numbers of the atom as a whole.
This physical picture allows us to conjecture that what is important for the slow precession rate are only the vectors $L$ and $S$ averaged over the period of the much faster precession about vector $J$ – in other words, only their components $L_J$ and $S_J$ along the vector $J$. Classically, these components may be calculated as

$$L_J = \frac{L \cdot J}{J^2} J, \quad \text{and} \quad S_J = \frac{S \cdot J}{J^2} J. \quad (6.71)$$

The scalar products participating in these expressions may be readily expressed via the squared lengths of the vectors, using the following evident classical formulas:

$$S^2 = (J - L)^2 \equiv J^2 + L^2 - 2L \cdot J, \quad L^2 = (J - S)^2 \equiv J^2 + S^2 - 2J \cdot S. \quad (6.72)$$

As a result, we get the following time average:

$$L_z + 2S_z = (L_J + 2S_J)_z = \left( \frac{L \cdot J}{J^2} J + 2 \frac{S \cdot J}{J^2} J \right)_z = \frac{J_z}{J} \left( \frac{L \cdot J}{J^2} J + 2 \frac{S \cdot J}{J^2} J \right)$$

$$= J_z \frac{(J^2 + L^2 - S^2) + 2(J^2 + S^2 - L^2)}{2J^2} \equiv J_z \left( 1 + \frac{J^2 + S^2 - L^2}{2J^2} \right). \quad (6.73)$$

The last move is to smuggle in some quantum mechanics by using, instead of the vector lengths squared, and the $z$-component of $J_z$, their eigenvalues given by Eqs. (5.169), (5.175), and (5.177). As a result, we immediately arrive at the exact result given by Eqs. (70). This coincidence encourages thinking about quantum mechanics of angular momenta in the classical terms of torque-induced precession, which turns out to be very fruitful in some more complex problems of atomic and molecular physics.

The high-field limit and low-field limits of the Zeeman effect, described respectively by Eqs. (66) and (69), are separated by a medium field range, in which the Zeeman splitting is of the order of the fine-structure splitting analyzed in Sec. 3. There is no time in this course for a quantitative analysis of this crossover.\textsuperscript{22}

6.5. Time-dependent perturbations

Now let us proceed to the case when the perturbation $\hat{H}^{(1)}$ in Eq. (1) is a function of time, while $\hat{H}^{(0)}$ is time-independent. The adequate perturbative approach to this problem, and its results, depend critically on the relation between the characteristic frequency (or the characteristic reciprocal time) $\omega$ of the perturbation and the distance between the initial system’s energy levels:

$$\hbar \omega \leftrightarrow |E_n - E_{n'}|. \quad (6.74)$$

In the case when all essential frequencies of a perturbation are very small in the sense of Eq. (74), we are dealing with the so-called adiabatic change of parameters, that may be treated essentially as a time-independent perturbation (see the previous sections of this chapter). The most interesting observation here is that the adiabatic perturbation does not allow any significant transfer of system’s

\textsuperscript{22} For a more complete discussion of the Stark, Zeeman, and fine-structure effects in atoms, I can recommend, for example, either the monograph by G. Woolgate cited above, or the one by I. Sobelman, Theory of Atomic Spectra, Alpha Science, 2006.
probability from one eigenstate to another. For example, in the WKB limit of the orbital motion, the Bohr quantization rule and its Wilson-Sommerfeld modification (2.110) guarantee that the integral
\[ \oint p \cdot dr , \] (6.75)
taken along the particle’s classical trajectory, is an adiabatic invariant, i.e. does not change at a slow change of system’s parameters. (It is curious that classical mechanics also guarantees the invariance of the integral (75), but its proof there\textsuperscript{23} is much harder than the quantum-mechanical derivation of this fact, carried out in Sec. 2.4.) This is why even if the perturbation becomes large with time (while changing sufficiently slowly), we can expect the eigenstate and eigenvalue classification to persist.

Let us proceed to the harder case when both sides of Eq. (74) are comparable, using for this discussion the Schrödinger picture of quantum dynamics, given by Eq. (4.158). Combining it with Eq. (1), we get the Schrödinger equation in the form
\[ i\hbar \frac{\partial}{\partial t} |\alpha(t)\rangle = \left( \hat{H}^{(0)} + H^{(1)}(t) \right) |\alpha(t)\rangle. \] (6.76)

Very much in the spirit of our treatment of the time-independent case in Sec. 1, let us represent the time-dependent ket-vector of the system with its expansion,
\[ |\alpha(t)\rangle = \sum_n |n\rangle \langle n| \alpha(t)\rangle , \] (6.77)
over the full and orthonormal set of the unperturbed, stationary ket-vectors defined by equation
\[ \hat{H}^{(0)} |n\rangle = E_n |n\rangle . \] (6.78)
(Note that these kets $|n\rangle$ are exactly what was called $|n^{(0)}\rangle$ in Sec. 1; we just may afford a less bulky notation in this section, because we will discuss only the lowest orders of the perturbation theory.) Plugging the expansion (77), with $n$ replaced with $n'$, into both sides of Eq. (76), and then inner-multiplying both its sides by the bra-vector $\langle n|$ of another unperturbed (and hence time-independent) state of the system, we get a set of linear, ordinary differential equations for the expansion coefficients:
\[ i\hbar \frac{d}{dt} \langle n| \alpha(t)\rangle = E_n \langle n| \alpha(t)\rangle + \sum_{n'} H_{nn'}^{(1)}(t) \langle n'| \alpha(t)\rangle, \] (6.79)
where the matrix elements of the perturbation in the unperturbed state basis, defined similarly to Eq. (8), are now functions of time:
\[ H_{nn'}^{(1)}(t) \equiv \langle n| \hat{H}^{(1)}(t) |n'\rangle . \] (6.80)

The set of differential equations (79), which are still exact, may be useful for numerical calculations.\textsuperscript{24} However, Eq. (79) has a certain technical inconvenience, which becomes clear if we consider its (evident) solution in the absence of perturbation:\textsuperscript{25}

\textsuperscript{23} See, e.g., CM Sec. 10.2.
\textsuperscript{24} Even if the problem under analysis may be described by the wave-mechanics Schrödinger equation (1.25), a direct numerical integration of that partial differential equation is typically less convenient than that of the ordinary differential equations (79).
\textsuperscript{25} This is of course just a more general form of Eq. (1.62) of the wave mechanics of time-independent systems.
\[ \langle n | \alpha(t) \rangle = \langle n | \alpha(0) \rangle \exp\left\{ -i \frac{E_n}{\hbar} t \right\}. \] (6.81)

We see that the solution oscillates very fast, and its numerical modeling may represent a challenge for even the fastest computers. These spurious oscillations (whose frequency, in particular, depends on the energy reference level) may be partly tamed by looking for the general solution of Eqs. (79) in a form inspired by Eq. (81):

\[ \langle n | \alpha(t) \rangle \equiv a_n(t) \exp\left\{ -i \frac{E_n}{\hbar} t \right\}. \] (6.82)

Here \( a_n(t) \) are new functions of time (essentially, the stationary states’ probability amplitudes), which may be used, in particular, to calculate the time-dependent level occupancies, i.e. the probabilities \( W_n \) to find the perturbed system on the corresponding energy levels of the unperturbed system:

\[ W_n(t) = \left| \langle n | \alpha(t) \rangle \right|^2 = \left| a_n(t) \right|^2. \] (6.83)

Plugging Eq. (82) into Eq. (79), for these functions we readily get a slightly modified system of equations:

\[ i\hbar \hat{a}_n = \sum_{n'} a_{n'} H_{nn'}^{(1)}(t) e^{i\omega_{nn'}t}, \] (6.84)

where the factors \( \omega_{nn'} \), defined by the relation

\[ \hbar \omega_{nn'} = E_n - E_{n'}, \] (6.85)

have the physical sense of frequencies of potential quantum transitions between the \( n \)-th and \( n' \)-th energy levels of the unperturbed system. (The conditions when such transitions indeed take place will be clear soon.) The advantages of Eq. (84) over Eq. (79), for both analytical and numerical calculations, is their independence of the energy reference, and lower frequencies of oscillations of the right-hand side terms, especially when the energy levels of interest are close to each other.\(^{26}\)

In order to continue our analytical treatment, let us restrict ourselves to a particular but very important case of a sinusoidal perturbation turned on at some moment – which may be taken for \( t = 0 \):

\[ \hat{H}^{(1)}(t) = \begin{cases} 0, & \text{for } t < 0, \\ \hat{A} e^{-i\omega t} + \hat{A}^\dagger e^{i\omega t}, & \text{for } t \geq 0, \end{cases} \] (6.86)

where the perturbation amplitude operators \( \hat{A} \) and \( \hat{A}^\dagger \),\(^{27}\) and hence their matrix elements,

---

\(^{26}\) Note that the relation of Eq. (84) to the initial Eq. (79) is very close to the relation of the interaction picture of quantum dynamics, discussed at the end of Sec. 4.6, to its Schrödinger picture, with the perturbation Hamiltonian playing the role of the interaction one – compare Eqs. (1) and Eq. (4.206). Indeed, Eq. (84) could be readily obtained from the interaction picture, and I did not do this just to avoid using this heavy bra-ket artillery for our (relatively) simple problem, and hence to keep its physics more transparent.

\(^{27}\) The notation of the amplitude operators in Eq. (86) is justified by the fact that the perturbation Hamiltonian has to be self-adjoint (Hermitian), and hence each term on the right-hand side of that relation has to be a Hermitian conjugate of its counterpart, which is evidently true only if the amplitude operators are also the Hermitian conjugates of each other. Note, however, that each of these amplitude operators is generally not Hermitian.
\[ \langle n | \hat{A} | n' \rangle \equiv A_{nn'}, \quad \langle n | \hat{A}^\dagger | n' \rangle = A_{n'n}^*, \quad (6.87) \]

are time-independent after the turn-on moment. In this case, Eq. (84) yields

\[ i\hbar \dot{a}_n = \sum_{n'} a_{n'n'} e^{i(\omega_{n'n} - \omega)t} + A_{n'n}^* e^{i(\omega_{n'n} + \omega)t}, \quad \text{for } t > 0. \quad (6.88) \]

This is, generally, still a nontrivial system of coupled differential equations; however, it allows simple and explicit solutions in two very important limits. First, let us assume that our system initially was definitely in one eigenstate \( n' \) (usually, though not necessarily, in the ground state), and that the occupancies \( W_n \) of all other levels stay very low all the time. (We will find the condition when the second assumption is valid \textit{a posteriori} – from the solution.) With the corresponding assumption

\[ a_{n'} = 1; \quad |a_n| << 1, \quad \text{for } n \neq n', \quad (6.89) \]

Eq. (88) may be readily integrated, giving

\[ a_n = -\frac{A_{nn'}}{\hbar(\omega_{nn'} - \omega)} \left[ e^{i(\omega_{nn'} - \omega)t} - 1 \right] - \frac{A_{n'n}}{\hbar(\omega_{n'n} + \omega)} \left[ e^{i(\omega_{n'n} + \omega)t} - 1 \right], \quad \text{for } n \neq n'. \quad (6.90) \]

This expression describes what is colloquially called the \textit{ac} excitation of (other) energy levels. Qualitatively, it shows that the probability \( W_n \) (83) of finding the system in each state (“on each energy level”) of the system does not tend to any constant value but rather oscillates in time – the so-called \textit{Rabi oscillations}. It also shows that that the \textit{ac}-field-induced transfer of the system from one state to the other one has a clearly resonant character: the maximum occupancy \( W_n \) of a level grows infinitely when the corresponding \textit{detuning}28

\[ \Delta_{nn'} \equiv \omega - \omega_{nn'}, \quad (6.91) \]

tends to zero. This conclusion is clearly unrealistic, and is an artifact of our initial assumption (89); according to Eq. (90), it is satisfied only if\(^{29}\)

\[ |A_{nn'}| << \hbar |\omega \pm \omega_{nn'}|, \quad (6.92) \]

and hence which does not allow a more deep analysis of the resonant excitation.

In order to overcome this limitation, we may perform the following trick – very similar to the one we used for the transfer to the degenerate case in Sec. 1. Let us assume that for a certain level \( n \),

\[ |\Delta_{nn'}| << \omega, |\omega \pm \omega_{nn'}|, |\omega \pm \omega_{n'n'}|, \quad \text{for all } n'' \neq n, n'. \quad (6.93) \]

- the condition illustrated in Fig. 8. Then, according to Eq. (90), we may ignore the occupancy of all but two levels, \( n \) and \( n' \), and also the second, non-resonant term with frequency \( \omega_{nn'} + \omega \approx 2\omega >> |\Delta_{nn'}| \) in Eqs. (88),\(^{30}\) now written for two probability amplitudes, \( a_n \) and \( a_{n'} \).

28 The notion of detuning is also very useful in the classical theory of oscillations (see, e.g., CM Chapter 5), where the role of \( \omega_{nn'} \) is played by the own frequency \( \omega_0 \) of the oscillator.

29 Strictly speaking, one more condition is that the number of “resonance” levels is also not too high – see Sec. 6.

30 The second assumption, i.e. the omission of non-resonant terms in the equations for the amplitudes is called the \textit{Rotating Wave Approximation} (RWA); the same idea in the classical theory of oscillations is the basis of what is usually called the \textit{van der Pol method}, and its result, the \textit{reduced equations} – see, e.g., CM Secs. 5.3-5.5.
The result is the following system of two linear equations:

\[ i\hbar \dot{a}_n = a_n A e^{-i\Delta t}, \quad i\hbar \dot{a}^*_{n'} = a^*_n A e^{i\Delta t}, \]  

which uses the shorthand notation \( A \equiv A_{nn'} \) and \( \Delta \equiv \Delta_{nn'} \). (I will use it for a while – until other energy levels become involved, at the beginning of the next section). This system may be readily reduced to a form without explicit time dependence of the right-hand parts – for example, by introducing the following new probability amplitudes, which the same moduli:

\[ b_n = a_n e^{i\Delta t / 2}, \quad b^*_{n'} = a^*_n e^{-i\Delta t / 2}, \]  

so that

\[ a_n = b_n e^{-i\Delta t / 2}, \quad a^*_{n'} = b^*_{n'} e^{i\Delta t / 2}. \]  

Plugging these relations into Eq. (94), we get two usual linear first-order differential equations:

\[ i\hbar \dot{b}_n = -\frac{\hbar\Delta}{2} b_n + A b^*_{n'}, \quad i\hbar \dot{b}^*_{n'} = A^* b_n + \frac{\hbar\Delta}{2} b^*_{n'}. \]  

As the reader knows very well by now, the general solution of such a system is a linear combination of two exponential functions, \( \exp\{\lambda_{\pm} t\} \), with the exponents \( \lambda_{\pm} \) that may be found by plugging any of these functions into Eq. (97), and requiring the consistency of two resulting linear algebraic equations. In our case, the consistency condition (i.e. the characteristic equation of the system) is

\[ \begin{vmatrix} -\frac{\hbar\Delta}{2} - i\hbar\lambda_{-} & A \\ A^* & \frac{\hbar\Delta}{2} - i\hbar\lambda_{+} \end{vmatrix} = 0, \]  

with roots \( \lambda_{\pm} = \pm i\Omega \), where

\[ \Omega \equiv \left( \frac{\Delta^2}{4} + \frac{|A|^2}{\hbar^2} \right)^{1/2}, \quad \text{i.e.} \quad 2\Omega = \left( \frac{\Delta^2}{4} + 4 \frac{|A|^2}{\hbar^2} \right)^{1/2}. \]  

The coefficients at the exponents are determined by initial conditions. If, as was assumed before, the system was completely on the level \( n' \) initially, at \( t = 0 \), i.e. \( a_{n'}(0) = 1, a_n(0) = 0 \), so that \( b_{n'}(0) = 1, b_n(0) = 0 \) as well, then Eqs. (97) yield, in particular:

\[ b_n(t) = -i \frac{A}{\hbar\Omega} \sin \Omega t, \]  

so that the \( n^{th} \) level occupancy is
\[ W_n = \left| h_n \right|^2 = \frac{|A|^2}{\hbar^2 \Omega^2} \sin^2 \Omega t \equiv \frac{|A|^2}{|A|^2 + (\hbar \Delta/2)^2} \sin^2 \Omega t. \] (6.101)

This is the famous *Rabi oscillation formula*.\(^{31}\) If the detuning is large in comparison with \( |A|/\hbar \), though still small in the sense of Eq. (93), the frequency \( 2\Omega \) of the Rabi oscillations is completely determined by the detuning, and their amplitude is small:

\[ W_n(t) = 4 \frac{|A|^2}{\hbar^2 \Delta^2} \sin^2 \frac{\Delta t}{2} \ll 1, \quad \text{for } |A|^2 \ll (\hbar \Delta)^2, \] (6.102)

- the result which could be obtained directly from Eq. (90), just neglecting the second term on its right-hand side. However, now we may also analyze the results of an increase of the perturbation amplitude \( |A| \): it leads not only to an increase of the amplitude of the probability oscillations, but also of their frequency – see Fig. 9. Ultimately, at \( |A| >> \hbar |\Delta| \) (for example, at the exact resonance, \( \Delta = 0 \), i.e. \( \omega_{n'} = \omega \), so that \( E_n = E_{n'} + \hbar \omega \)), Eqs. (101)-(102) give \( \Omega = |A|/\hbar \) and \( (W_n)_{\text{max}} = 1 \), i.e. describe a periodic, full “repumping” of the system from one level to another and back, with a frequency proportional to the perturbation amplitude.\(^{32}\)

This effect is a close analog of the quantum oscillations in two-level systems with time-independent Hamiltonians, which were discussed in Secs. 2.6 and 5.1. Indeed, let us revisit, for a moment, their discussion started at the end of Sec.1, now paying more attention to the time evolution of the system under the perturbation. As was argued in that section, the most general perturbation Hamiltonian lifting the two-fold degeneracy of an energy level, in an arbitrary basis, has the matrix (28). Let us describe the system’s dynamics using, again, the Schrödinger picture, representing the ket-vector of an arbitrary state of the system in the form (5.1), where \( \uparrow \) and \( \downarrow \) are the time-independent states of the

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\(^{31}\) It was derived in 1952 by Isaac Rabi, in the context of his group’s pioneering experiments with the ac (practically, microwave) excitation of quantum states, using molecular beams in vacuum.

\(^{32}\) As Eqs. (82), (96), and (99) show, the lowest frequency in the system is \( \omega_k = \omega_{n'} - \Delta/2 + \Omega \), so that at \( A \to 0 \), \( \hbar \omega_k \approx \hbar \omega_{n'} + 2 |A|^2/\hbar \Delta \). This effective shift of the lowest energy level (which may be measured by another “probe” field of a different frequency) is a particular case of the *ac Stark effect*, which was already mentioned in Sec. 2.
basis in which Eq. (28) is written (now without any obligation to associate these states with the \(z\)-basis of any spin-\(\frac{1}{2}\)). Then, the Schrödinger equation (4.158) yields

\[
\begin{bmatrix}
H(\hat{\alpha}_\uparrow)
\end{bmatrix} = H(\hat{\alpha}_\downarrow)
\end{bmatrix} = \begin{bmatrix}
H_{11} & H_{12} \\
H_{21} & H_{22}
\end{bmatrix}
\begin{bmatrix}
\alpha_\uparrow \\
\alpha_\downarrow
\end{bmatrix} = \begin{bmatrix}
H_{11}\alpha_\uparrow + H_{12}\alpha_\downarrow \\
H_{21}\alpha_\uparrow + H_{22}\alpha_\downarrow
\end{bmatrix}.
\]

(6.103)

As we know (for example, from the discussion in Sec. 5.1), the average of the diagonal elements of the matrix gives just a common shift of the system’s energy; for the purpose of the dynamics analysis, it may be absorbed into the energy reference level. Also, the Hamiltonian operator has to be Hermitian, so that the off-diagonal elements of its matrix have to be complex-conjugate. With this, Eqs. (103) are reduced to the form,

\[
\begin{align*}
\ih \dot{\alpha}_\uparrow &= -\frac{\xi}{2}\alpha_\uparrow + H_{12}\alpha_\downarrow, \\
\ih \dot{\alpha}_\downarrow &= H*_{12}\alpha_\uparrow + \frac{\xi}{2}\alpha_\downarrow, \\
\end{align*}
\]

with \(\hbar\xi \equiv H_{22} - H_{11},\)

(6.104)

which is absolutely similar to Eqs. (97). In particular, these equations describe the quantum oscillations of the probabilities \(W_\uparrow = |\alpha_\uparrow|^2\) and \(W_\downarrow = |\alpha_\downarrow|^2\) with the frequency\(^{33}\)

\[
2\Omega = \left(\frac{\xi^2 + 4|H_{12}|^2}{\hbar^2}\right)^{1/2}.
\]

(6.105)

The similarity of Eqs. (97) and (104), and hence of Eqs. (99) and (105), shows that the Rabi oscillations and the “usual” quantum oscillations have essentially the same physical nature, besides that in the former case the external ac signal quantum \(\hbar\omega\) bridges the separated energy levels, effectively reducing their difference \((E_n - E_{n'})\) to a much smaller difference \(-\Delta \equiv (E_n - E_{n'}) - \hbar\omega\). Also, since the Hamiltonian (28) is similar to that given by Eq. (5.2), the dynamics of such a system with two ac-coupled energy levels, within the limits (93) of the perturbation theory, is completely similar to that of a time-independent two-level system. In particular, its state may be similarly represented by a point on the Bloch sphere shown in Fig. 5.3, with its dynamics described, in the Heisenberg picture, by Eq. (5.19). This fact is very convenient for the experimental implementation of quantum information systems (to be discussed in more detail in Sec. 8.5), because it enables qubit manipulations in a broad variety of physical systems with well-separated energy levels, using external ac (usually microwave or optical) sources.

Note, however, that according to Eq. (90), if the system has energy levels other than \(n\) and \(n'\), they also become occupied to some extent. Since the sum of all occupancies equals 1, this means that \((W_n)_{\text{max}}\) may approach 1 only if the other excitation amplitude is very small, and hence the state manipulation time scale \(\tau = 2\pi/\Omega = 2\pi\hbar/|A|\) is very long. The ultimate limit in this sense is provided by the harmonic oscillator where all energy levels are equidistant, and the probability repumping between all of them occurs at an approximately the same rate. In particular, in this system the implementation of the full Rabi oscillations is impossible even at the exact resonance.\(^{34}\)

\(^{33}\) By the way, Eq. (105) gives a natural generalization of the relations obtained for the frequency of such oscillations in Sec. 2.6, where the coupled potential wells were assumed to be exactly similar, so that \(\xi = 0\). Moreover, Eqs. (104) gives a long-promised proof of Eqs. (2.201), and hence a better justification of Eqs. (2.203).

\(^{34}\) From Sec. 5.5, we already know what happens to the ground state of an oscillator at its external sinusoidal (or any other) excitation: it turns into a Glauber state, i.e. a superposition of all Fock states – see Eq. (5.134).
However, I would not like these quantitative details to obscure from the reader the most important qualitative (OK, maybe semi-quantitative :-) conclusion of this section’s analysis: a resonant increase of the interlevel transition intensity at \( \omega \rightarrow \omega_{nn'} \). As will be shown later in the course, in a quantum system coupled to its environment at least slightly (hence in reality, in any quantum system), such increase is accompanied by a sharp increase of the external field’s absorption, which may be measured. This effect has numerous practical applications including systems based on the electron paramagnetic resonance (EPR) and nuclear magnetic resonance (NMR) spectroscopies, which are broadly used in material science, chemistry, and medicine. Unfortunately, I will not have time to discuss the related technical issues and methods (in particular, interesting ac pulsing techniques, including the so-called Ramsey interferometry) in detail, and have to refer the reader to special literature.35

6.6. Quantum-mechanical Golden Rule

One of the results of the past section, Eq. (102), may be used to derive one of the most important and nontrivial results of quantum mechanics. For that, let us consider the case when the perturbation causes quantum transitions from a discrete energy level \( E_{n'} \) into a group of eigenstates with a dense (essentially continuous) spectrum \( E_n \) – see Fig. 10a.

If, for all states \( n \) of the group, the following conditions are satisfied

\[
|A_{nn'}|^2 \ll (\hbar \Delta_{nn'})^2 \ll (\hbar \omega_{nn'})^2,
\]

then Eq. (102) coincides with the result that would follow from Eq. (90). This means that we may apply Eq. (102), with the indices \( n \) and \( n' \) duly restored, to any level \( n \) of our tight group. As a result, the total probability of having our system transferred from the initial level \( n' \) to that group is

\[
W_z(t) = \sum_n W_n(t) = \frac{4}{\hbar^2} \sum_n \frac{|A_{nn'}|^2}{\Delta_{nn'}^2} \sin^2 \frac{\Delta_{nn'} t}{2}.
\]

Now comes the main, absolutely beautiful trick: let us assume that the summation over \( n \) is limited to a tight group of very similar states whose matrix elements \( A_{nn'} \) are virtually similar (we will check the validity of this assumption later on), so that we can take it out of the sum (107) and then replace the sum with the corresponding integral:

---

\[ W_{n'}(t) = \frac{4|A_{nn'}|^2}{\hbar^2} \int \frac{1}{\Delta_{nn'}^2} \sin^2 \frac{\Delta_{nn'} t}{2} d\Delta_{nn'} = \frac{4|A_{nn'}|^2}{\hbar} \rho_{n'} t \int \frac{1}{(\Delta_{nn'}^2)^2} \sin^2 \frac{\Delta_{nn'} t}{2} \, d(-\Delta_{nn'} t), \quad (6.108) \]

where \( \rho_{n'} \) is the density of the states on the energy axis:

\[ \rho_{n'} = \frac{dn}{dE_n}. \quad (6.109) \]

This density and the matrix element \( A_{nn'} \) have to be evaluated at \( \Delta_{nn'} = 0 \), i.e. at energy \( E_n = E_{n'} + \hbar \omega \), and are assumed to be constant within the final state group. At fixed \( E_{n'} \), the function under integral (108) is even and decreases fast at \( |\Delta_{nn'} t| \gg 1 \) – see Fig. 10b. Hence we may introduce a dimensionless integration variable \( \xi \equiv \Delta_{nn'} t \), and extend the integration over it formally from \(-\infty\) to \(+\infty\). Then the integral in Eq. (108) is reduced to a table one, and yields

\[ W_{n'}(t) = \frac{4|A_{nn'}|^2}{\hbar} \rho_{n'} t + \int_{-\infty}^{\infty} \frac{1}{\xi^2} \sin^2 \frac{\xi}{2} \, d\xi = \frac{4|A_{nn'}|^2}{\hbar} \rho_{n'} \frac{\pi}{2} = \Gamma t, \quad (6.110) \]

where the constant

\[ \Gamma = \frac{2\pi}{\hbar} |A_{nn'}|^2 \rho_{n'} \quad (6.111) \]

is called the transition rate. This is one of the most famous and useful results of quantum mechanics, its Golden Rule, which deserves much discussion.

First of all, let us reproduce the reasoning already used in Sec. 2.5 to show that the meaning of the rate \( \Gamma \) is much deeper than Eq. (110) seems to imply. Indeed, due to the conservation of the total probability, \( W_{n'} + W_{n} = 1 \), we can rewrite that equation as

\[ \dot{W}_{n'} \big|_{t=0} = -\Gamma. \quad (6.112) \]

Evidently, this result cannot be true for all times, otherwise the probability \( W_{n'} \) would become negative. The reason for this apparent contradiction is that Eq. (110) was obtained in the assumption that initially, the system was completely on level \( n' \): \( W_{n'}(0) = 1 \). Now, if at the initial moment the value of \( W_{n'} \) is different, the result (110) has to be multiplied by that number, due to the linear relation (88) between \( da_n/dt \) and \( a_n' \). Hence, instead of Eq. (112) we get a differential equation similar to Eq. (2.159),

\[ \dot{W}_{n'} \big|_{t=0} = -\Gamma W_{n'}, \quad (6.113) \]

which, for a time-independent \( \Gamma \), has the evident solution,

---

36 See, e.g., MA Eq. (6.12).

37 In some texts, the density of states in Eq. (111) is replaced with a formal expression \( \sum \delta(E_n - E_{n'} - \hbar \omega) \). Indeed, applied to a finite energy interval \( \Delta E_n \) with \( \Delta n >> 1 \) levels, it gives the same result: \( \Delta n \equiv (dn/dE_n)\Delta E_n \equiv \rho_{n}\Delta E_n \). Such replacement may be technically useful in some cases, but is incorrect for \( \Delta n \sim 1 \), and hence should be used with the utmost care, so that for most applications the more explicit form (111) is preferable.

38 Sometimes Eq. (111) is called “Fermi’s Golden Rule”. This is rather unfair, because this result was developed mostly by the same P. A. M. Dirac in 1927, and Enrico Fermi’s role was not much more than advertising it, under the name of “Golden Rule No. 2”, in his influential lecture notes on nuclear physics, which were published much later, in 1950. (To be fair to Fermi, he has never tried to pose as the Golden Rule’s author.)
describing the exponential decay of the initial state’s occupancy, with the time constant \( \tau = 1/\Gamma \).

I am inviting the reader to review this fascinating result again: by the summation of periodic oscillations (102) over many levels \( n \), we have got an exponential decay (114) of the probability. This trick becomes possible because the effective range \( \Delta E_n \) of the state energies \( E_n \) giving substantial contributions into the integral (108), shrinks with time: \( \Delta E_n \sim \hbar/t \). By the way, since most of the decay takes place within the time interval of the order of \( \tau \equiv 1/\Gamma \), the range of the participating final energies may be estimated as

\[
\Delta E_n \sim \frac{\hbar}{\tau} \equiv \hbar \Gamma .
\]

This estimate is very instrumental for the formulation of conditions of the Golden Rule’s validity. First, we have assumed that the matrix elements of the perturbation and the density of states are independent of the energy within the interval (115). This gives the following requirement

\[
\Delta E_n \sim \hbar \Gamma < \ll E_n - E_{n'} \sim \hbar \omega,
\]

Second, for the transfer from the sum (107) to the integral (108), we need the number of states within that energy interval, \( \Delta N_n = \rho_n \Delta E_n \), to be much larger than 1. Merging Eq. (116) with Eq. (92) for all the energy levels \( n'' \neq n, n' \) not participating in the resonant transition, we may summarize all conditions of the Golden Rule validity as

\[
\rho_n^{-1} \ll \hbar \Gamma \ll \hbar |\omega \pm \omega_{n''}| .
\]

(The reader may ask whether I have neglected the condition expressed by the first of Eqs. (106). However, for \( \Delta_{n''} \sim \Delta E_{n''} / \hbar \sim \Gamma \), this condition is just \( |A_{nn''}|^2 < (\hbar \Gamma)^2 \), so that plugging it into Eq. (111),

\[
\Gamma < \frac{2\pi}{\hbar} (\hbar \Gamma)^2 \rho_n ,
\]

and canceling one \( \Gamma \) and one \( \hbar \), we see that it coincides with the first relation in Eq. (117) above.)

Let us have a look at whether these conditions may be satisfied in practice, at least in some cases. For example, let us consider the optical ionization of an atom, with the released electron confined in a volume of the order of 1 cm\(^3\) \( \equiv 10^{-6} \) m\(^3\). According to Eq. (1.90), with \( E \) of the order of the atomic ionization energy \( E_n - E_{n'} = \hbar \omega \sim 1 \) eV, the density of electron states in that volume is of the order of \( 10^{21} \) 1/eV, while the right-hand side of Eq. (117) is of the order of \( E_n \sim 1 \) eV. Thus the conditions (117) provide an approximately 20-orders-of-magnitude range for acceptable values of \( \hbar \Gamma \). This illustration should give the reader a taste of why the Golden Rule is applicable to so many situations.

Finally, the physical picture of the initial state’s decay (which will also be the key to our discussion of quantum-mechanical “open” systems in the next chapter) is also very important. According to Eq. (114), the external excitation transfers the system into the continuous spectrum of levels \( n \), and it never comes back to the initial level \( n' \). However, it was derived from the quantum mechanics of the initial state's decay...
mechanics of Hamiltonian systems, whose equations are invariant with respect to time reversal. This paradox is a result of the generalization (113) of the exact result (112), which breaks the time-reversal symmetry. This is a trick of course, but the tick absolutely adequate for the physics under study. Indeed, some gut feeling of the physical sense of this irreversibility may be obtained from the following observation. As Eq. (1.86) illustrates, the distance between the adjacent orbital energy levels tends to zero only if the system’s size goes to infinity. This means that our assumption of the continuous energy spectrum of the final states \( n \) essentially requires these states to be broadly extended in space – being either free, or essentially free de Broglie waves. Thus the Golden Rule corresponds to the (physically justified) assumption that in an infinitely large system, the traveling de Broglie waves excited by a local source and propagating outward from it, would never come back, and even if they did, the unpredictable phase shifts introduced by the uncontrollable perturbations on their way would never allow them to sum up in the coherent way necessary to bring the system back into the initial state \( n' \). (This is essentially the same situation which was discussed, for a particular 1D wave-mechanical system, in Sec. 2.5.)

To get a feeling of the Golden Rule at work, let us apply it to the following simple problem – which is a toy model of the photoelectric effect, briefly discussed in Sec. 1.1(ii). A 1D particle is initially trapped in the ground state of a narrow potential well, described by Eq. (2.158):

\[
U(x) = -\omega \delta(x), \quad \text{with } \omega > 0.
\]  

Let us calculate the rate \( \Gamma \) of the particle’s “ionization” (i.e. its excitation into a group of extended, delocalized states) by a weak classical sinusoidal force of amplitude \( F_0 \) and frequency \( \omega \), suddenly turned on at some instant.

As a reminder, the initial, localized state (in our current notation, \( n' \)) of such a particle was already found in Sec. 2.6:

\[
\psi_n(x) = \kappa^{1/2} \exp\left(-\kappa |x|\right), \quad \text{with } \kappa \equiv \frac{m \omega}{\hbar^2}, \quad E_{n'} \equiv -\frac{\hbar^2 \kappa^2}{2m} = -\frac{m \omega^2}{2\hbar^2}.
\]  

The final, extended states \( n \), with a continuous spectrum, for this problem exist only at energies \( E_n > 0 \), so that the excitation rate is different from zero only for frequencies

\[
\omega > \omega_{\text{min}} \equiv \frac{|E_{n'}|}{\hbar} = \frac{m \omega^2}{2\hbar^3}.
\]  

The weak sinusoidal force may be described by the following perturbation Hamiltonian,

\[
\hat{H}^{(1)} = -F(t) \hat{x} = -F_0 \dot{x} \cos \omega t \equiv -\frac{F_0}{2} \dot{x} \left(e^{i\omega t} + e^{-i\omega t}\right), \quad \text{for } t > 0,
\]  

so that according to Eq. (86), which serves as the amplitude operator’s definition, in this case

\[
\hat{A} = \hat{A}^\dagger = -\frac{F_0}{2} \hat{x}.
\]  

Now the matrix elements \( A_{nn'} \) that participate in Eq. (111) may be calculated in the coordinate representation:

\[40\] This situation is also very much similar to the entropy increase in macroscopic systems, which is postulated in thermodynamics, and justified in statistical physics, even though it is based on time-reversible laws of mechanics – see, e.g., SM Secs. 1.2 and 2.2.
\[
A_{nm'} = \int_{-\infty}^{+\infty} \psi_n^*(x) \hat{A}(x) \psi_{n'}(x) dx = -\frac{F_0}{2} \int_{-\infty}^{+\infty} \psi_n^*(x) x \psi_{n'}(x) dx. \quad (6.124)
\]

Since, according to Eq. (120), the initial \( \psi_n \) is a symmetric function of \( x \), non-vanishing contributions to this integral are given only by asymmetric functions \( \psi_{n'}(x) \), proportional to \( \sin k_n x \), with the wave number \( k_n \) related to the final energy by the well-familiar equality (1.89):

\[
\frac{\hbar^2 k_n^2}{2m} = E_n. \quad (6.125)
\]

As we know from Sec. 2.6 (see in particular Eq. (2.167) and its discussion), such asymmetric functions, with \( \psi_n(0) = 0 \), are not affected by the zero-centered delta-functional potential (119), so that their density \( \rho_n \) is the same as that in completely free space, and we can use Eq. (1.100). (Actually, since that relation was derived for traveling waves, it is more prudent to repeat the calculation that has led to that result, confining the waves on an artificial segment \([-l/2, +l/2]\) – so long, with

\[
k_n l, \kappa l \gg 1, \quad (6.126)
\]

that it does not affect the initial localized state and the excitation process. Then the confinement requirement \( \psi_n(\pm l/2) = 0 \) immediately yields the condition \( k_n l/2 = n\pi \), so that Eq. (1.100) is indeed valid, but only for positive values of \( k_n \) because \( \sin k_n x \) with \( k_n \to -k_n \) does not describe an independent standing-wave eigenstate.) Hence the final state density is

\[
\rho_n \equiv \frac{dn}{dE_n} = \frac{dn}{dk_n} \left/ \frac{dE_n}{dk_n} \right/ \frac{h^2 k_n}{2m} \equiv \frac{lm}{2\pi\hbar^2 k_n}. \quad (6.127)
\]

It may look troubling that the density of states depends on the artificial segment’s length \( l \), but the same \( l \) also participates in the final wavefunctions’ normalization factor,41

\[
\psi_n = \left( \frac{2}{l} \right)^{1/2} \sin k_n x, \quad (6.128)
\]

and hence in the matrix element (124):

\[
A_{nm'} = -\frac{F_0}{2} \left( \frac{2\kappa}{l} \right)^{1/2} \int_{-l}^{l} \sin k_n x e^{-\kappa|x|} x dx = -\frac{F_0}{2i} \left( \frac{2\kappa}{l} \right)^{1/2} \left( \int_{0}^{l} e^{(i\kappa - \kappa)x} x dx - \int_{0}^{-l} e^{-(i\kappa + \kappa)x} x dx \right). \quad (6.129)
\]

These two integrals may be readily worked out by parts. Taking into account that, according to the condition (126), their upper limits may be extended to \( \infty \), the result is

\[
A_{nm'} = -\left( \frac{2\kappa}{l} \right)^{1/2} \frac{2k_n \kappa}{(k_n^2 + \kappa^2)^2}. \quad (6.130)
\]

Note that the matrix element is a smooth function of \( k_n \) (and hence of \( E_n \)), so that the main assumption of the Golden Rule, the virtual constancy of \( A_{nn'} \) in the interval \( \Delta E_n \sim \hbar \Gamma \ll E_{n'} \), is satisfied. So, Eq. (111) is reduced to the following expression:

---

41 The normalization to infinite volume, using Eq. (4.263), is also possible, but physically less transparent.
\[
\Gamma = \frac{2\pi}{\hbar} \left( \frac{2\kappa}{l} \right)^{1/2} \left( \frac{2k_n\kappa}{(k_n^2 + \kappa^2)^2} \right) \frac{2m}{\hbar^2 k_n} \equiv \frac{8F_0^2 mk_n\kappa^3}{h^3(k_n^2 + \kappa^2)^2}, \tag{6.131}
\]

which is independent of the artificially introduced \(l\) – thus justifying its use.

Note that due to the above definitions of \(k_n\) and \(\kappa\), the expression in parentheses in the denominator of the last expression does not depend on the potential well’s “weight” \(\nu\), and is a function of only the excitation frequency \(\omega\) (and the particle’s mass):

\[
\frac{\hbar^2(k_n^2 + \kappa^2)}{2m} = E_n - E_{n'} = \hbar\omega. \tag{6.132}
\]

As a result, Eq. (131) may be recast simply as

\[
\hbar\Gamma = \frac{F_0^2 \omega^3 k_n}{2(h\omega)^4}. \tag{6.133}
\]

What is hidden here is that \(k_n\), defined by Eq. (125) with \(E_n = E_{n'} + \hbar\omega\), is a function of the external force’s frequency, changing as \(\omega^{1/2}\) at \(\omega >> \omega_{\text{min}}\) (so that \(\Gamma\) drops as \(\omega^{-7/2}\) at \(\omega \to \infty\)), and as \((\omega - \omega_{\text{min}})^{1/2}\) when \(\omega\) approaches the “red boundary” (121) of the ionization effect, so that \(\Gamma \propto (\omega - \omega_{\text{min}})^{1/2} \to 0\) in that limit as well. So, our toy model does describe this main feature of the photoelectric effect, whose explanation by A. Einstein was essentially the starting point of quantum mechanics – see Sec. 1.1.

The conceptually very similar, but a bit more involved analysis of such effect in a more realistic 3D case, namely the hydrogen atom’s ionization by an optical wave, is left for the reader’s exercise.

### 6.7. Golden Rule for step-like perturbations

Now let us reuse some of our results for a perturbation being turned on at \(t = 0\), but after that time-independent:

\[
\hat{H}^{(1)}(t) = \begin{cases} 
0, & \text{for } t < 0, \\
\hat{H} = \text{const}, & \text{for } t \geq 0.
\end{cases} \tag{6.134}
\]

A superficial comparison of this equality and the former Eq. (86) seems to indicate that we may use all our previous results, taking \(\omega = 0\) and replacing \(\hat{A} + \hat{A}^\dagger\) with \(\hat{H}^{(1)}\). However, that conclusion (which would give us a wrong factor of 2 in the result) does not take into account the fact that analyzing both the two-level approximation in Sec. 5, and the Golden Rule in Sec. 6, we have dropped the second (non-resonant) term in Eq. (90). This why it is more prudent to use the general Eq. (84),

\[
i\hbar\hat{a}_n = \sum_{n'} a_{nn'}H_{\omega_m}e^{i\omega_{m'}t}, \tag{6.135}
\]

in which the matrix element of the perturbation is now time-independent at \(t > 0\). We see that it is formally equivalent to Eq. (88) with only the first (resonant) term kept, if we make the following replacements:

\[
\hat{A} \to \hat{H}, \quad \Delta_{mm'} \equiv \omega - \omega_{mm'} \to -\omega_{mm'}. \tag{6.136}
\]
Let us use this equivalency to consider the results of coupling between a discrete-energy state \( n' \), into which the particle is initially placed, and a dense group of states with a quasi-continuum spectrum, in the same energy range. Figure 11a shows an example of such a system: a particle is initially (say, at \( t = 0 \)) placed into a potential well separated by a penetrable potential barrier from a formally infinite region with a continuous energy spectrum. Let me hope that the physical discussion in the last section makes the outcome of such an experiment evident: the particle will gradually and irreversibly tunnel out of the well, so that the probability \( W_{n'}(t) \) of its still residing in the well will decay in accordance with Eq. (114). The rate of this decay may be found by making the replacements (136) in Eq. (111):

\[
\Gamma = \frac{2\pi}{\hbar} |H_{nn'}|^2 \rho_n,
\]

where the states \( n \) and \( n' \) now have virtually the same energy.\(^{42}\)

It is very informative to compare this result, semi-quantitatively, with Eq. (105) for a symmetric \( (E_n = E_{n'}) \) system of two potential wells separated by a similar potential barrier – see Fig. 11b. For the symmetric case, i.e. \( \xi = 0 \), Eq. (105) is reduced to simply

\[
\Omega = \frac{1}{\hbar} |H_{nn'}|_{\text{con}}.
\]

Here I have used the index “con” (from “confinement”) to emphasize that this matrix element is somewhat different from the one participating in Eq. (137), even if the potential barriers are similar. Indeed, in the latter case, the matrix element,

\[
H_{nn'} = \langle n | \hat{H} | n' \rangle = \int \psi_n^* \hat{H} \psi_{n'} dx,
\]

has to be calculated for two wavefunctions \( \psi_n \) and \( \psi_{n'} \), confined to spatial intervals of the same scale \( l_{\text{con}} \), while in Eq. (137), the wavefunctions \( \psi_n \) are extended over a much larger distance \( l >> l_{\text{con}} \) – see Fig. 11. As Eq. (128) tells us, in the 1D model this means an additional small factor of the order of \( (l_{\text{con}}/l)^{1/2} \). Now using Eq. (128) as a crude but suitable model for the final-state wavefunctions, we arrive at the following estimate, independent of the artificially introduced length \( l \):

\[
\hbar \Gamma \sim 2\pi |H_{nn'}|_{\text{con}}^2 \frac{l_{\text{con}}}{l} \rho_n \sim 2\pi |H_{nn'}|_{\text{con}}^2 \frac{l_{\text{con}}}{l} \frac{\ln m}{2\pi \hbar^2 k_n} \frac{|H_{nn'}|^2_{\text{con}}}{\Delta E_{n'}} \equiv (\hbar \Omega)^2, \quad (6.140)
\]

where \( \Delta E_{n'} \sim \hbar^2 / m l_{\text{con}}^2 \) is the scale of the differences between the eigenenergies of the particle in an unperturbed potential well. Since the condition of validity of Eq. (138) is \( \hbar \Omega << \Delta E_{n'} \), we see that

\[^{42}\text{The condition of validity of Eq. (137) is again given by Eq. (117), just with } \omega = 0 \text{ in the upper limit for } \Gamma.\]
This (sufficiently general\textsuperscript{43}) perturbative result confirms the conclusion of a more particular analysis carried out in the end of Sec. 2.6: the rate of the (irreversible) quantum tunneling into a state continuum is always much lower than the frequency of (reversible) quantum oscillations between discrete states separated with the same potential barrier – at least for the case when both are much lower than $\Delta E_n / \hbar$, so that the perturbation theory is valid. A very handwaving interpretation of this result is that the particle oscillates between the confined state in the well and the space-extended states behind the barrier many times before finally “deciding to perform” an irreversible transition into the unconfined continuum. This qualitative picture is consistent with experimentally observable effects of dispersive electromagnetic environments on electron tunneling.\textsuperscript{44}

Let me conclude this section (and this chapter) with the application of Eq. (137) to a very important case, which will provide a smooth transition to the next chapter’s topic. Consider a composite system consisting of two “component” systems, $a$ and $b$, with the energy spectra sketched in Fig. 12.

![Fig. 6.12. Energy relaxation in system $a$ due to its weak coupling to system $b$ (which serves as the environment of $a$).](image)

Let the systems be completely independent initially. The independence means that in the absence of their coupling, the total Hamiltonian of the system may be represented as a sum of two operators:

$$\hat{H}^{(0)} = \hat{H}_a(a) + \hat{H}_b(b),$$

where the arguments $a$ and $b$ symbolize the non-overlapping sets of the degrees of freedom of the two systems. Such operators, belonging to their individual, different Hilbert spaces, naturally commute. Similarly, the eigenkets of the system may be naturally factored as

$$|n\rangle = |n_a\rangle \otimes |n_b\rangle.$$  

(6.143)

The direct product sign $\otimes$ is used here (and below) to denote the formation of a joint ket-vector from the kets of the independent systems, belonging to different Hilbert spaces. Evidently, the order of operands in such a product may be changed at will. As a result, its eigenenergies separate into a sum, just as the Hamiltonian (142) does:

$$\hat{H}^{(0)}|n\rangle = \left(\hat{H}_a + \hat{H}_b\right)|n_a\rangle \otimes |n_b\rangle \equiv \left(\hat{H}_a|n_a\rangle\right) \otimes |n_b\rangle + \left(\hat{H}_b|n_b\rangle\right) \otimes |n_a\rangle = \left(E_{na} + E_{nb}\right)|n\rangle.$$  

(6.144)

\textsuperscript{43} It is straightforward to verify that the estimate (141) is valid for similar problems of any spatial dimensionality, not just for the 1D case we have analyzed.

In such composite systems, the relatively weak interaction of its components may be usually represented as a bilinear product of two Hermitian operators, each depending only on the degrees of freedom of only one component system:

$$\hat{H}^{(1)} = \hat{A}(a)\hat{B}(b).$$

(6.145)

A very common example of such an interaction is the electric-dipole interaction between an atomic-scale system (with a linear size of the order of the Bohr radius $r_B \sim 10^{-10}$ m) and the electromagnetic field at optical frequencies $\omega \sim 10^{16}$ s$^{-1}$, with the wavelength $\lambda = 2\pi/\omega \sim 10^{-6}$ m >> $r_B$.\(^{45}\)

$$\hat{H}^{(1)} = -\hat{d} \cdot \hat{\mathbf{e}}, \quad \text{with} \quad \hat{d} \equiv \sum_k q_k \hat{\mathbf{r}}_k,$$

(6.146)

where the dipole electric moment $\mathbf{d}$ depends only on the positions $\mathbf{r}_k$ of the charged particles (numbered with index $k$) of the atomic system, while that of electric field $\mathbf{e}$ is a function of only the electromagnetic field’s degrees of freedom – to be discussed in Chapter 9 below.

Returning to the general situation shown in Fig. 12, if the component system $a$ was initially in an excited state $n'a$, the interaction (145), turned on at some moment of time, may bring it into another discrete state $n_a$ of a lower energy – for example, the ground state. In the process of this transition, the released energy, in the form of an energy quantum

$$\hbar \omega \equiv E_{n'a} - E_{na},$$

(6.147)

is picked up by the system $b$:

$$E_{n'b} = E_{n'a} + \hbar \omega \equiv E_{n'a} + (E_{n'a} - E_{na}),$$

(6.148)

so that the total energy $E = E_a + E_b$ of the system does not change. (If the states $n_a$ and $n'b$ are the ground states of the two component systems, as they are in most applications of this analysis, and we take the ground state energy $E_g = E_{n'a} + E_{n'b}$ of the composite system for the reference, then Eq. (148) gives merely $E_{n'b} = E_{n'a}$.) If the final state $n_b$ of the system $b$ is inside a state group with a quasi-continuous energy spectrum (Fig. 12), the process has the exponential character (114)\(^{46}\) and may be interpreted as the effect of energy relaxation of the system $a$, with the released energy quantum $\hbar \omega$ absorbed by the system $b$. Note that since the quasi-continuous spectrum essentially requires a system of a large spatial size, such a model is very convenient for description of the environment $b$ of the quantum system $a$. (In physics, the “environment” typically means all the Universe – less the system under consideration.)

If the relaxation rate $\Gamma$ is sufficiently low, it may be described by the Golden Rule (137). Since the perturbation (145) does not depend on time explicitly, and the total energy $E$ does not change, this relation, with the account of Eqs. (143) and (145), takes the form

$$\Gamma = \frac{2\pi}{\hbar} |A_{nn'}|^2 |B_{nn'}|^2 \rho_n, \quad \text{where} \quad A_{nn'} = \langle n_a | \hat{A} | n'a \rangle, \quad \text{and} \quad B_{nn'} = \langle n_b | \hat{B} | n'b \rangle,$$

(6.149)

where $\rho_n$ is the density of the final states of the system $b$ at the relevant energy (147). In particular, Eq. (149), with the dipole Hamiltonian (146), enables a very straightforward calculation of the natural linewidth of atomic electric-dipole transitions. However, such calculation has to be postponed until Chapter 9, in which we will discuss the electromagnetic field quantization – i.e., the exact nature of the

\(^{45}\) See, e.g., EM Sec. 3.1, in particular Eq. (3.16), in which letter $p$ is used for the electric dipole moment.

\(^{46}\) Such process is spontaneous: it does not require any external agent, and starts as soon as either the interaction (145) has been turned on, or (if it is always on) as soon as the system $a$ is placed into the excited state $n'a$.\(\)
states $n_b$ and $n'_b$ for this problem, and hence will be able to calculate $B_{nn'}$ and $\rho_n$. Instead, I will now proceed to a general discussion of the effects of quantum systems interaction with their environment, toward which the situation shown in Fig. 12 provides a clear conceptual path.

### 6.8. Exercise problems

6.1. Use Eq. (14) to prove the following general form of the Hellmann-Feynman theorem (whose proof in the wave-mechanics domain was the task of Problem 1.5):

$$\frac{\partial E_n}{\partial \lambda} = \langle n | \frac{\partial \hat{H}}{\partial \lambda} | n \rangle,$$

where $\lambda$ is an arbitrary $c$-number parameter.

6.2. Establish a relation between Eq. (16) and the result of the classical theory of weakly anharmonic (“nonlinear”) oscillations at negligible damping.

*Hint:* Use N. Bohr’s reasoning discussed in Problem 1.1.

6.3. A weak, time-independent force $F$ is exerted on a 1D particle that was placed into a hard-wall potential well

$$U(x) = \begin{cases} 0, & \text{for } 0 < x < a, \\ +\infty, & \text{otherwise.} \end{cases}$$

Calculate, sketch, and discuss the 1st-order perturbation of its ground-state wavefunction.

6.4. A time-independent force $F = \mu(n_x + n_y)$, where $\mu$ is a small constant, is applied to a 3D harmonic oscillator of mass $m$ and frequency $\omega_0$. Calculate, in the first order of the perturbation theory, the effect of the force upon the ground state energy of the oscillator, and its lowest excited energy level. How small should the constant $\mu$ be for your results to be quantitatively correct?

6.5. A 1D particle of mass $m$ is localized at a narrow potential well that may be approximated with a delta-function:

$$U(x) = -\omega \delta(x), \quad \text{with } \omega > 0.$$ 

Calculate the change of its ground state’s energy by an additional weak, time-independent force $F$, in the first non-vanishing approximation of the perturbation theory. Discuss the limits of validity of this result, taking into account that at $F \neq 0$, the localized state of the particle is metastable.

6.6. Use the perturbation theory to calculate the eigenvalues of the operator $\hat{L}^2$ in the limit $|m| \approx l >> 1$, by purely wave-mechanical means.

*Hint:* Try the following substitution: $\Theta(\theta) = f(\theta)/\sin^{1/2} \theta$.

6.7. In the first non-vanishing order of the perturbation theory, calculate the shift of the ground-state energy of an electrically charged spherical rotator (i.e. a particle of mass $m$, free to move over a spherical surface of radius $R$) due to a weak, uniform, time-independent electric field $\mathcal{E}$. 

6.8. Use the perturbation theory to evaluate the effect of a time-independent, uniform electric field \( \mathcal{E} \) on the ground state energy \( E_g \) of a hydrogen atom. In particular:

(i) calculate the 2\(^{nd}\)-order shift of \( E_g \), neglecting the extended unperturbed states with \( E > 0 \), and bring the result to the simplest analytical form you can,
(ii) find the lower and the upper bounds on the shift, and
(iii) discuss the simplest experimental manifestations of this quadratic Stark effect.

6.9. A particle of mass \( m \), with electric charge \( q \), is in its ground \( s \)-state with a given energy \( E_g < 0 \), being localized by a very short-range, spherically-symmetric potential well. Calculate its static electric polarizability \( \alpha \).

6.10. In some atoms, the charge-screening effect of other electrons on the motion of each of them may be reasonably well approximated by the replacement of the Coulomb potential (3.190), \( U = -C/r \), with the so-called Hulthén potential

\[
U = -\frac{C/a}{\exp[r/a]-1} \rightarrow -C \times \begin{cases} 1/r, & \text{for } r << a, \\ \exp[-r/a]/a, & \text{for } a << r. 
\end{cases}
\]

Assuming that the effective screening radius \( a \) is much larger than \( r_0 \equiv \hbar^2/mC \), use the perturbation theory to calculate the energy spectrum of a single particle of mass \( m \), moving in this potential, in the lowest order needed to lift the \( l \)-degeneracy of the levels.

6.11. In the first non-vanishing order of the perturbation theory, calculate the correction to energies of the ground state and all lowest excited states of a hydrogen-like atom/ion, due to electron’s penetration into its nucleus, modeling it as a spinless, uniformly charged sphere of a radius \( R << r_B/Z \).

6.12. Prove that the kinetic-relativistic correction operator (48) indeed has only diagonal matrix elements in the basis of unperturbed Bohr atom states (3.200).

6.13. Calculate the lowest-order relativistic correction to the ground-state energy of a 1D harmonic oscillator.

6.14. Use the perturbation theory to calculate the contribution to the magnetic susceptibility \( \chi_m \) of a dilute gas, that is due to the orbital motion of a single electron inside each gas particle. Spell out your result for a spherically-symmetric ground state of the electron, and give an estimate of the magnitude of this orbital susceptibility.

6.15. How to calculate the energy level degeneracy lifting, by a time-independent perturbation, in the 2\(^{nd}\) order of the perturbation in \( \hat{H}^{(1)} \), assuming that it is not lifted in the 1\(^{st}\) order? Carry out such calculation for a plane rotator of mass \( m \) and radius \( R \), carrying electric charge \( q \), and placed into a weak, uniform, constant electric field \( \mathcal{E} \).

6.16. The Hamiltonian of a quantum system is slowly changed in time.

(i) Develop a theory of quantum transitions in the system, and spell out its result in the 1\(^{st}\) order in the speed of the change.
(ii) Use the 1st-order result to calculate the probability that a finite-time pulse of a slowly changing force $F(t)$ drives a 1D harmonic oscillator, initially in its ground state, into an excited state.

(iii) Compare the last result with the exact one.

6.17. Use the single-particle model to calculate the complex electric permittivity $\varepsilon(\omega)$ of a dilute gas of similar atoms, due to their induced electric polarization by a weak external ac field, for a field frequency $\omega$ very close to one of quantum transition frequencies $\omega_{n'n}$. Based on the result, calculate and estimate the absorption cross-section of each atom.

*Hint:* In the single-particle model, atom’s properties are determined by $Z$ similar, non-interacting electrons, each moving in a similar static attracting potential, generally different from the Coulomb one, because it is contributed not only by the nucleus, but also by other electrons.

6.18. Use the solution of the previous problem to generalize the expression for the London dispersion force between two atoms (whose calculation in the harmonic-oscillator model was the subject of Problems 3.16 and 5.15) to the single-particle model with an arbitrary energy spectrum.

6.19. Use the solution of the previous problem to calculate the potential energy of interaction of two hydrogen atoms, both in their ground state, separated by distance $r >> r_B$.

6.20. In a certain quantum system, distances between the three lowest energy levels are slightly different – see the figure on the right (| $\xi$ | << $\omega_{1,2}$). Assuming that the involved matrix elements of the perturbation Hamiltonian are known, and are all proportional to the external ac field’s amplitude, find the time necessary to populate the first excited level almost completely (with a given precision $\varepsilon << 1$), using the Rabi oscillation effect, if at $t = 0$ the system is completely in its ground state.

6.21. Analyze the possibility of a slow transfer of a system from one of its energy levels to another one (in the figure on the right, from level 1 to level 3), using the scheme shown in that figure, in which the monochromatic external excitation amplitudes $A_+$ and $A_-$ may be slowly changed at will.

6.22. A weak external force pulse $F(t)$, of a finite time duration, is applied to a 1D harmonic oscillator that initially was in its ground state.

(i) Calculate, in the lowest non-vanishing order of the perturbation theory, the probability that the pulse drives the oscillator into its lowest excited state.

(ii) Compare the result with the exact solution of the problem.

(iii) Spell out the perturbative result for a Gaussian-shaped waveform,

$$F(t) = F_0 \exp\left(-t^2 / \tau^2\right),$$

and analyze its dependence on the scale $\tau$ of the pulse duration.

6.23. A spatially-uniform, but time-dependent external electric field $E(t)$ is applied, starting from $t = 0$, to a charged plane rotator, initially in its ground state.
(i) Calculate, in the lowest non-vanishing order in the field’s strength, the probability that by 
time $t > 0$, the rotator is in its $n^{th}$ excited state.

(ii) Spell out and analyze your results for a constant-magnitude field rotating, with a constant 
angular velocity $\omega$, within the rotator’s plane.

(iii) Do the same for a monochromatic field of frequency $\omega$, with a fixed polarization.

6.24. A spin-$\frac{1}{2}$ with a gyromagnetic ratio $\gamma$ is placed into a magnetic field including a time-
independent component $B_0$, and a perpendicular field of a constant magnitude $B_r$, rotated with a 
constant angular velocity $\omega$. Can this magnetic resonance problem be reduced to one already discussed 
in Chapter 6?

6.25. Develop the general theory of quantum excitations of the higher levels of a discrete-
spectrum system, initially in the ground state, by a weak time-dependent perturbation, up to the 2nd 
order. Spell out and discuss the result for the case of a monochromatic excitation, with a nearly perfect 
tuning of its frequency $\omega$ to the half of a certain quantum transition frequency $\omega_0 \equiv (E_n - E_0)/\hbar$.

6.26. A heavy, relativistic particle, with electric charge $q = Ze$, passes by a hydrogen atom, 
initially in its ground state, with an impact parameter (the shortest distance) $b$ within the range $r_B << b 
<< r_B/\alpha$, where $\alpha \approx 1/137$ is the fine structure constant. Calculate the probabilities of the atom’s 
transition to its lowest excited states.

6.27. A particle of mass $m$ is initially in the localized ground state, with the known energy $E_g < 0$, of a very small, spherically-symmetric potential well. Calculate the rate of its delocalization by an 
applied classical force $F(t) = nF_0\cos \omega t$ with a time-independent direction $n$.

6.28.* Calculate the rate of ionization of a hydrogen atom, initially in its ground state, by a 
classical, linearly polarized electromagnetic wave with an electric field’s amplitude $\vec{E}_0$, and a frequency 
$\omega$ within the range

$$\frac{\hbar}{m_e r_B^2} << \omega << \frac{c}{r_B},$$

where $r_B$ is the Bohr radius. Recast your result in terms of the cross-section of this electromagnetic wave 
absorption process. Discuss briefly what changes of the theory would be necessary if either of the above 
conditions had been violated.

6.29.* Use the quantum-mechanical Golden Rule to derive the general expression for the electric 
current $I$ through a weak tunnel junction between two conductors, biased with dc voltage $V$, treating the 
conductors as Fermi gases of electrons, with negligible direct interaction. Simplify the result in the low-
voltage limit.

*Hint:* The electric current flowing through a weak tunnel junction is so low that it does not substantially perturb the electron states inside each conductor.

6.30.* Generalize the result of the previous problem to the case when a weak tunnel junction is 
biased with voltage $V(t) = V_0 + A\cos \omega t$, with $\hbar \omega$ generally comparable with $eV_0$ and $eA$.

6.31.* Use the quantum-mechanical Golden Rule to derive the Landau-Zener formula (2.257).