## 4. Angular momentum

In the following we consider general and exact approaches, such as angular momentum algebra and group theory, and later, mehods for approximations, such as perturbation theory, in applications of quantum mechanics.
Consider first the angular momentum (liikemäärämomentti, kulmaliikemäärä, pyörimismäärä) and derivation of its quantization from commutation relation, only.

## Angular momentum operators

These operators can be written in terms of position and momentum operators q and $\mathrm{p}_{\mathrm{q}}$, whose commutation relations are

$$
\begin{equation*}
\left[\mathrm{q}, \mathrm{p}_{\mathrm{q}^{\prime}}\right]=\mathrm{i} \hbar \delta_{\mathrm{qq}}{ }^{\prime} \tag{4.1}
\end{equation*}
$$

where $\mathrm{q}, \mathrm{q}^{\prime}=\mathrm{x}, \mathrm{y}$ or z .

### 4.1. Operators and their commutation relations

Classically angular momentun is defined as

$$
\ell=\mathbf{r} \times \mathbf{p}=\left|\begin{array}{ccc}
\hat{\mathbf{i}} & \hat{\mathbf{j}} & \hat{\mathbf{k}}  \tag{4.2}\\
\mathrm{x} & \mathrm{y} & \mathrm{z} \\
p_{x} & p_{y} & p_{z}
\end{array}\right|=\left(y p_{z}-z p_{y}\right) \hat{\mathbf{i}}+\left(z p_{x}-x p_{z}\right) \hat{\mathbf{j}}+\left(x p_{y}-y p_{x}\right) \hat{\mathbf{k}},
$$

and thus, its components are

$$
\begin{equation*}
\ell_{\mathrm{x}}=\mathrm{yp}_{\mathrm{z}}-\mathrm{zp} p_{\mathrm{y}}, \quad \ell_{\mathrm{y}}=\mathrm{zp} p_{\mathrm{x}}-\mathrm{xp} p_{\mathrm{z}} \text { and } \ell_{\mathrm{z}}=\mathrm{xp} p_{\mathrm{y}}-\mathrm{yp} p_{\mathrm{x}} \tag{4.3}
\end{equation*}
$$

Classically $\ell=\mathbf{I} \cdot \omega$ ja $\left.\mathrm{E}=\ell^{2} / 2 \mathrm{I}=\frac{1}{2} \mathrm{I} \omega^{2}\right)$ and

$$
\begin{equation*}
\ell^{2}=\ell_{\mathrm{x}}^{2}+\ell_{\mathrm{y}}^{2}+\ell_{\mathrm{z}}^{2} \tag{4.4}
\end{equation*}
$$

Now, as in quantum mechanics

$$
\left(p_{x} \rightarrow \frac{\hbar}{\mathrm{i}} \frac{\partial}{\partial \mathrm{x}}\right), \quad\left(\mathrm{p}_{\mathrm{y}} \rightarrow \frac{\hbar}{\mathrm{i}} \frac{\partial}{\partial \mathrm{y}}\right) \quad \text { and } \quad\left(\mathrm{p}_{\mathrm{z}} \rightarrow \frac{\hbar}{\mathrm{i}} \frac{\partial}{\partial \mathrm{z}}\right),
$$

it follows that the angular momentum components take the form

$$
\begin{align*}
& \ell_{\mathrm{x}}=\frac{\hbar}{\mathrm{i}}\left(\mathrm{y} \frac{\partial}{\partial \mathrm{z}}-\mathrm{z} \frac{\partial}{\partial \mathrm{y}}\right), \\
& \ell_{\mathrm{y}}=\frac{\hbar}{\mathrm{i}}\left(\mathrm{z} \frac{\partial}{\partial \mathrm{x}}-\mathrm{x} \frac{\partial}{\partial \mathrm{z}}\right) \text { and }  \tag{4.5}\\
& \ell_{\mathrm{z}}=\frac{\hbar}{\mathrm{i}}\left(\mathrm{x} \frac{\partial}{\partial \mathrm{y}}-\mathrm{y} \frac{\partial}{\partial \mathrm{x}}\right) .
\end{align*}
$$

Let us find commutators:
[ $\ell_{\mathrm{x}}, \ell_{\mathrm{y}}$ ] =

Similarly, the other two, and thus,

$$
\begin{equation*}
\left[\ell_{\mathrm{x}}, \ell_{\mathrm{y}}\right]=\mathrm{i} \hbar \ell_{z},\left[\ell_{\mathrm{y}}, \ell_{\mathrm{z}}\right]=\mathrm{i} \hbar \ell_{\mathrm{x}} \mathrm{ja}\left[\ell_{z}, \ell_{\mathrm{x}}\right]=\mathrm{i} \hbar \ell_{\mathrm{y}}, \tag{4.7}
\end{equation*}
$$

which are the basic commutation rules of angular momentum in quantum mechanics.
As above, it can be shown that

$$
\begin{equation*}
\left[\ell^{2}, \ell_{\mathrm{q}}\right]=0 . \tag{4.8}
\end{equation*}
$$

### 4.2. Angular momentum observables

In a compact form (4.7) can be written as

$$
\begin{equation*}
\ell \times \ell=\mathrm{i} \hbar \ell, \tag{4.9}
\end{equation*}
$$

because

$$
\ell \times \ell=\left|\begin{array}{ccc}
\mathbf{i} & \mathbf{j} & \mathbf{k} \\
\ell_{\mathrm{x}} & \ell_{\mathrm{y}} & \ell_{\mathrm{z}} \\
\ell_{\mathrm{x}} & \ell_{\mathrm{y}} & \ell_{\mathrm{z}}
\end{array}\right|=
$$

### 4.3. Shift operators

Let us define raising and lowering operator, which are also called

with inverse relations

$$
\begin{equation*}
\ell_{+}=\ell_{\mathrm{x}}+\mathrm{i} \ell_{\mathrm{y}} \quad \text { and } \quad \ell_{-}=\ell_{\mathrm{x}}-\mathrm{i} \ell_{\mathrm{y}} \tag{4.10}
\end{equation*}
$$

It is easy to show that

$$
\left[\ell_{+}, \ell_{2}\right]=-\hbar \ell_{+},\left[\ell_{-}, \ell_{2}\right]=-\hbar \ell_{-} \text {ja }\left[\ell_{+}, \ell_{-}\right]=2 \hbar \ell_{z},
$$

and that $\ell^{2}$ commutes with ladder operators

$$
\begin{equation*}
\left[\ell^{2}, \ell_{ \pm}\right]=0 . \tag{4.12}
\end{equation*}
$$

Note that shift operators are not hermitian, but complex conjugates of each other.

## Definition of states

Next, we will find the angular momentum states by using the commutation relations.

### 4.4. Effect of shift operators

Because $\ell^{2}$ and $\ell_{\mathrm{z}}$ commute, (4.8), they have the same eigenfunctions. Thus, from (3.33) and (3.37) we have the eigenvalue equations

$$
\begin{align*}
& \ell^{\ell^{2}\left|\ell, \mathrm{~m}_{\ell}\right\rangle=\hbar^{2} \ell(\ell+1)\left|\ell, \mathrm{m}_{\ell}\right\rangle \text { and }}  \tag{4.14-15}\\
& \ell_{\mathrm{z}}\left|\ell, \mathrm{~m}_{\ell}\right\rangle=\hbar \mathrm{m}_{\ell}\left|\ell, \mathrm{m}_{\ell}\right\rangle,
\end{align*}
$$

where $\ell=0,1,2,3, \ldots$ and $\mathrm{m}_{\ell}=\ell, \ell-1, \ell-2, \ldots,-\ell$. Let us find next $\ell_{+}\left|\ell, \mathrm{m}_{\ell}\right\rangle$. From (4.12)

$$
\begin{aligned}
\ell_{+} \ell_{\mathrm{z}}-\ell_{\mathrm{z}} \ell_{+} & =-\hbar \ell_{+} \text {, thus } \ell_{\mathrm{z}} \ell_{+}\left|\ell_{,} \mathrm{m}_{\ell}\right\rangle=\left(\ell_{+} \ell_{\mathrm{z}}+\hbar \ell_{+}\right)\left|\ell, \mathrm{m}_{\ell}\right\rangle \\
& =\ell_{+} \hbar \mathrm{m}_{\ell}\left|\ell, \mathrm{m}_{\ell}\right\rangle+\hbar \ell_{+}\left|\ell, \mathrm{m}_{\ell}\right\rangle=\hbar\left(\mathrm{m}_{\ell}+1\right) \ell_{+}\left|\ell, \mathrm{m}_{\ell}\right\rangle .
\end{aligned}
$$

Therefore, $\ell_{+}\left|\ell_{\ell} \mathrm{m}_{\ell}\right\rangle$ is an eigenfunction of $\ell_{\mathrm{z}}$ with the eigenvalue $\hbar\left(\mathrm{m}_{\ell}+1\right)$. So,
and similarly

$$
\begin{align*}
& \ell_{+}\left|\ell, \mathrm{m}_{\ell}\right\rangle=\text { constant } \times\left|\ell, \mathrm{m}_{\ell}+1\right\rangle  \tag{4.17a}\\
& \ell-\left|\ell, \mathrm{m}_{\ell}\right\rangle=\text { constant } \times\left|\ell, \mathrm{m}_{\ell}-1\right\rangle \tag{4.17b}
\end{align*}
$$

These equations hold, if $\mathrm{m}_{\ell}<\ell$ and $\mathrm{m}_{\ell}>-\ell$, and for $\ell_{+}|\ell, \ell\rangle=0$ and $\ell_{-}\left|\ell_{1}-\ell\right\rangle=0$. For this reason $\ell_{+}$and $\ell_{-}$are called raising and lowering operators.

### 4.5. Eigenvalues of angular momentum

As discussed in sec. 3, operators $\ell^{2}$ and $\ell_{\mathrm{z}}$ correspond the "rotational states" of electrons on atomic orbitals, when $\ell=0,1,2,3, \ldots$. Then, quantization follows from the uniqueness of the wavefunction.

More generally, starting from the commutation relations (4.7), only, one can choose to accept half-integral quantum numbers, too. Then, the angular momentum can assume "half-integral" values, as well. For example, the spin of an electron or proton are half-integral angular momenta.

Let us denote the common eigenfunctions of general angular momentum j and its z component with $\left|\mathrm{j}, \mathrm{m}_{\mathrm{j}}\right\rangle$. Then, we can write the eigenvalue equations as

$$
\begin{equation*}
\mathrm{j}^{2}\left|\mathrm{j}, \mathrm{~m}_{\mathrm{j}}\right\rangle=\hbar^{2} \mathrm{j}(\mathrm{j}+1)\left|\mathrm{j}, \mathrm{~m}_{\mathrm{j}}\right\rangle ; \mathrm{j}=0, \frac{1}{2}, 1, \frac{3}{2}, 2, \ldots \text { and } \tag{4.22}
\end{equation*}
$$

### 4.6. Matrix elements of shift operators

The matrix elements of angular momentum can be obtained by using relations (4.10) and (4.11). Therefore, let us find $j_{ \pm}\left|j, m_{j}\right\rangle=c_{ \pm}\left|j, m_{j} \pm 1\right\rangle$. So,

$$
\mathrm{j}_{-} \mathrm{j}_{+}=
$$

Now, $\left(\mathrm{j}_{+}\left|\mathrm{j}, \mathrm{m}_{\mathrm{j}}\right\rangle\right)^{*}=\left\langle\mathrm{j}, \mathrm{m}_{\mathrm{j}}\right| \mathrm{j}_{-}$, and thus,

Similarly, we can derive

$$
\begin{align*}
& \left.\mathrm{j}^{+}\left|\mathrm{j}, \mathrm{~m}_{\mathrm{j}}\right\rangle=\hbar\left(\mathrm{j}(\mathrm{j}+1)-\mathrm{m}_{\mathrm{j}}\left(\mathrm{~m}_{\mathrm{j}}+1\right)\right)^{1 / 2} \mathrm{lj}, \mathrm{~m}_{\mathrm{j}}+1\right\rangle  \tag{4.29a}\\
& \left.\left.\mathrm{j}^{-} \mathrm{l}, \mathrm{~m}_{\mathrm{j}}\right\rangle=\hbar\left(\mathrm{j}(\mathrm{j}+1)-\mathrm{m}_{\mathrm{i}}\left(\mathrm{~m}_{\mathrm{j}}-1\right)\right)^{1 / 2} \mathrm{lj}, \mathrm{~m}_{\mathrm{i}}-1\right\rangle
\end{align*}
$$

### 4.7. Orbital angular momentum eigenfunctions

As an example, let us consider orbital angular momentum $\ell$, whose eigenfunctions are the spherical harmonics, in sec. 3 these results were found directly from the Schrödinger equation. Now, let us find these results from the general properties of angular momentum. In polar coordinates the angular momentum operators are

$$
\begin{align*}
\ell_{\mathrm{x}} & =-\frac{\hbar}{\mathrm{i}}\left(\sin \phi \frac{\partial}{\partial \theta}+\cot \theta \cos \phi \frac{\partial}{\partial \phi}\right) \\
\ell_{\mathrm{y}} & =\frac{\hbar}{\mathrm{i}}\left(\cos \phi \frac{\partial}{\partial \theta}-\cot \theta \sin \phi \frac{\partial}{\partial \phi}\right)  \tag{4.30}\\
\ell_{\mathrm{z}} & =\frac{\hbar}{\mathrm{i}} \frac{\partial}{\partial \phi}
\end{align*}
$$

and

$$
\begin{equation*}
\ell^{+}=\ell_{\mathrm{x}}+\mathrm{i} \ell_{\mathrm{y}}=\hbar \mathrm{e}^{\mathrm{i} \phi}\left(\frac{\partial}{\partial \theta}+\mathrm{i} \cot \theta \frac{\partial}{\partial \phi}\right) . \tag{4.31}
\end{equation*}
$$

Thus, for the eigenfunctions of the state state $\mathrm{m}_{\ell}=\ell$ we can write $\ell_{+}|\ell, \ell\rangle=0$ or $\ell^{+}=\ell_{\mathrm{x}}+\mathrm{i} \ell_{\mathrm{y}}=\hbar \mathrm{e} \mathrm{i}\left(\frac{\partial}{\partial \theta}+\mathrm{i} \cot \theta \frac{\partial}{\partial \phi}\right)$. whose solutions are the spherical harmonics, for $\mathrm{m}_{\ell}=\ell$ because

$$
\begin{equation*}
\psi_{\ell \ell}(\theta, \phi)=\mathrm{Y}_{\ell \ell}(\theta, \phi)=\mathrm{N} \sin ^{\ell} \theta \mathrm{e}^{\mathrm{i} \ell \phi} \tag{4.32}
\end{equation*}
$$

Now, the eigenfunctions for the other values of $\mathrm{m}_{\ell}$ are easily obtained by using $\ell_{-}$.

### 4.8. Spin

In 1925 Uhlenbeck and Goudsmit suggested an internal angular momentum of the electron, its spin, and its only quantum number $1 / 2$, to simplify interpretation of atomic spectra. Later, Dirac "found" the half-integer quantum numbers and the electron spin, in particular, from his relativistic extension of quantum mechanics.
So, the electron spin is $1 / 2$ and denoted as $s=1 / 2$. Therefore, its z-component is $\mathrm{m}_{\mathrm{s}}= \pm 1 / 2$ (up or down). The corresponding eigenstates are $\left|\mathrm{ls}, \mathrm{m}_{\mathrm{s}}\right\rangle$, usually written as

$$
\alpha=\left|\frac{1}{2}, \frac{1}{2}\right\rangle \text { and } \beta=\frac{1}{2}\left|-\frac{1}{2}\right\rangle \text {. }
$$

Eigenvalue equations are

$$
\begin{array}{|c|}
\hline \mathrm{s}_{\mathrm{z}} \alpha=\frac{1}{2} \hbar \alpha, \quad \mathrm{~s}_{\mathrm{z}} \beta=-\frac{1}{2} \hbar \beta,  \tag{4.33}\\
\mathrm{~s}^{2} \alpha=\frac{3}{4} \hbar^{2} \alpha \text { and } \mathrm{s}^{2} \beta=\frac{3}{4} \hbar^{2} \beta \\
\hline
\end{array}
$$

and by using the ladder operators, we obtain

$$
\begin{align*}
& \mathrm{s}^{+} \alpha=0, \quad \mathrm{~s}^{+} \beta=\hbar \alpha, \\
& \mathrm{s}^{-} \alpha=\hbar \beta \text { and } \mathrm{s}^{-} \beta=0 \tag{4.34}
\end{align*}
$$

and thus, the non-zero matrix elements are

$$
\begin{equation*}
\langle\alpha| s+|\beta\rangle=\hbar \text { and }\langle\beta| s-|\alpha\rangle=\hbar . \tag{4.35}
\end{equation*}
$$

## Coupling of angular momenta in composite systems

Nex, consider a system formed by two angular momenta, e.g. orbital angular momenta of two electrons $\ell$, or the orbital angular momentum $\ell$ and spin s of a single electron.

### 4.9. Uncoupled and coupled states

The state of two angular momenta can be defined by "listing the quantum numbers" in ket vector $\left.I \mathrm{j}_{1}, \mathrm{~m}_{\mathrm{j} 1} ; \mathrm{j}_{2}, \mathrm{~m}_{\mathrm{j} 2}\right\rangle$. It can be done so, because

$$
\mathrm{j}_{1}^{2}, \mathrm{j}_{1 \mathrm{z}}, \mathrm{j}_{2}^{2} \text { and } \mathrm{j}_{2 \mathrm{z}} \text { all commute. }
$$

For finding the total angular momentum "vector" j , the two vectors $\mathrm{j}_{1}$ and $\mathrm{j}_{2}$ are "summed componentwise", and the sum can be written as $\left.I \mathrm{j}_{1}, \mathrm{j}_{2} ; \mathrm{j}, \mathrm{m}_{\mathrm{j}}\right\rangle$, because
$\mathrm{j}_{1}{ }^{2}, \mathrm{j}_{2}{ }^{2}, \mathrm{j}^{2}$ and $\mathrm{j}_{\mathrm{z}}$ all commute. This is called coupling of $\mathrm{j}_{1}$ and $\mathrm{j}_{2}$.
However, generally

$$
\begin{equation*}
\left[j_{1 z}, j^{2}\right] \neq 0 \text { and }\left[j_{2 z}, j^{2}\right] \neq 0 \tag{4.37}
\end{equation*}
$$

and therefore, $\mathrm{m}_{\mathrm{j} 1}$ and $\mathrm{m}_{\mathrm{j} 2}$ can not be used to describe the state simultaneously with j . Later we will see, that depending on the case it is better to use the uncoupled representation $\left.\ \mathrm{j}_{1}, \mathrm{~m}_{\mathrm{j} 1} ; \mathrm{j}_{2}, \mathrm{~m}_{\mathrm{j} 2}\right\rangle$ or the coupled representation $\left.\mathrm{I} \mathrm{j}_{1}, \mathrm{j}_{2} ; \mathrm{j}, \mathrm{m}_{\mathrm{j}}\right\rangle$.
Exercise: Find out, if $\mathrm{j}=\mathrm{j}_{1}+\mathrm{j}_{2}$ is an angular momentum?

### 4.10. Permitted values of total angular momentum

Let us find the permitted values of j and $\mathrm{m}_{\mathrm{j}}$. Because

$$
\begin{aligned}
\mathrm{j}_{\mathrm{z}}\left|\mathrm{j}_{1}, \mathrm{~m}_{\mathrm{j} 1} ; \mathrm{j}_{2}, \mathrm{~m}_{\mathrm{j} 2}\right\rangle & =\left(\mathrm{j}_{\mathrm{z} 1}+\mathrm{j}_{22}\right)\left|\mathrm{j}_{1}, \mathrm{~m}_{\mathrm{j} 1} ; \mathrm{j}_{2}, \mathrm{~m}_{\mathrm{j} 2}\right\rangle \\
& =\hbar\left(\mathrm{m}_{\mathrm{j} 1}+\mathrm{m}_{\mathrm{j} 2}\right)\left|\mathrm{j}_{1}, \mathrm{~m}_{\mathrm{j} 1} ; \mathrm{j}_{2}, \mathrm{~m}_{\mathrm{j} 2}\right\rangle
\end{aligned}
$$

and $\left.\mathrm{j}_{\mathrm{z}}\left|\mathrm{j}_{1}, \mathrm{j}_{2} ; \mathrm{j}, \mathrm{m}_{\mathrm{j}}\right\rangle=\hbar \mathrm{m}_{\mathrm{j}} \backslash \mathrm{j}_{1}, \mathrm{j}_{2} ; \mathrm{j}, \mathrm{m}_{\mathrm{j}}\right\rangle$, we get

$$
\begin{equation*}
\mathrm{m}_{\mathrm{j}}=\mathrm{m}_{\mathrm{j} 1}+\mathrm{m}_{\mathrm{j} 2} . \tag{4.41}
\end{equation*}
$$

As the maximum value of $m_{j}$ is $j_{1}+j_{2}, j$ is allowed to assume values from "Clebsh-Gordan series"

$$
\begin{equation*}
\mathrm{j}=\mathrm{j}_{1}+\mathrm{j}_{2}, \mathrm{j}_{1}+\mathrm{j}_{2}-1, \ldots,\left|\mathrm{j}_{1}-\mathrm{j}_{2}\right| . \tag{4.42}
\end{equation*}
$$

The lower limit of Clebsh-Gordan series $\left|\mathrm{j}_{1}-\mathrm{j}_{2}\right|$ is defined by the condition, that the coupled presentation has to contain the same number of states as the uncoupled one does. This can be formulated as sc. "triangle condition".

Example Coupling of orbital angular momentum and spin of a p-electron $\ell=1$ and $\mathrm{s}=1 / 2$ in hydrogen atom.

### 4.11. Vector model of coupled angular momenta

Coupling of two angular momenta can also be graphically illustrated. It should fullfil the following conditions:

1. $|\mathbf{j}|=\sqrt{ }(\mathrm{j}(\mathrm{j}+1))$, where j is one of the permitted values in Clebsh-Gordan series.
2. For the $\mathbf{j}$-vector only the z -component is given, but not the x and $\mathrm{y} \Rightarrow$ cone description.
3. For $\mathbf{j}_{1}$ and $\mathbf{j}_{2}$ :
$\left|\mathbf{j}_{1}\right|=\sqrt{ }\left(\mathrm{j}_{1}\left(\mathrm{j}_{1}+1\right)\right)$ and
$\left|\mathbf{j}_{2}\right|=\sqrt{ }\left(\mathrm{j}_{2}\left(\mathrm{j}_{2}+1\right)\right)$
$\Rightarrow$ cone description.
4. The z-components $\mathrm{m}_{\mathrm{j} 1}$ and $\mathrm{m}_{\mathrm{j} 2}$ of $\mathbf{j}_{1}$ and $\mathbf{j}_{2}$ are shown in uncoupled representation.
a)

(b)



The z-component $m_{j}$ of $\mathbf{j}$ is given in the coupled representation.

Figs. 4.6-7.

(b)

(c)


Next, let us couple two spins $s_{1}=1 / 2$ and $s_{2}=1 / 2$. Uncoupled presentation consists of four possible states $\left.1 \mathrm{~s}_{1} \mathrm{~m}_{\mathrm{s} 1} ; \mathrm{s}_{2} \mathrm{~m}_{\mathrm{s} 2}\right\rangle$ as

$$
\begin{aligned}
& \left|\frac{1}{2} \frac{1}{2} ; \frac{1}{2} \frac{1}{2}\right\rangle=\alpha_{1} \alpha_{2} \quad\left|\frac{1}{2} \frac{1}{2} ; \frac{1}{2}-\frac{1}{2}\right\rangle=\alpha_{1} \beta_{2} \\
& \left|\frac{1}{2}-\frac{1}{2} ; \frac{1}{2} \frac{1}{2}\right\rangle=\beta_{1} \alpha_{2} \quad\left|\frac{1}{2}-\frac{1}{2} ; \frac{1}{2}-\frac{1}{2}\right\rangle=\beta_{1} \beta_{2}
\end{aligned}
$$

In coupling, the total angular momentum becomes $S=s_{1}+s_{2}, \ldots,\left|s_{1}-s_{2}\right|=1,0$ and the state $S=1$ is called triplet, because $\mathrm{M}_{\mathrm{S}}=1,0,-1$; and the state $\mathrm{S}=0$ and $\mathrm{M}_{\mathrm{S}}=0$ is called singlet, correspondingly. Thus, the couple representation also has four states $\left|\mathrm{s}_{1} \mathrm{~s}_{1} ; \mathrm{S}_{\mathrm{S}}\right\rangle \equiv\left|\frac{1}{2} \frac{1}{2} ; \mathrm{S}_{\mathrm{S}}\right\rangle \equiv\left|\mathrm{S}_{\mathrm{S}}\right\rangle$ as

$$
|1,1\rangle
$$

$|0,0\rangle|1,0\rangle$


Fig. 4.8.

$|1,-1\rangle$


Fig. 4.9-10.

QTES, 2021

### 4.12. Clebsh-Gordan coefficients

The coupled state wavefunction $\left.I \mathrm{j}_{1} \mathrm{j}_{2} ; \mathrm{j} \mathrm{m}_{\mathrm{j}}\right\rangle$ can be presented with those of uncoupled ones $\left.1 \mathrm{j}_{1} \mathrm{~m}_{\mathrm{j} 1} ; \mathrm{j}_{2} \mathrm{~m}_{\mathrm{j} 2}\right\rangle$

$$
\begin{equation*}
\left|\mathrm{j}_{1} \mathrm{j}_{2} ; \mathrm{j} \mathrm{~m}_{\mathrm{j}}\right\rangle=\sum_{\substack{\mathrm{m}_{\mathrm{j}} \mathrm{~m}_{\mathrm{j} 2} \\ \mathrm{~m}_{\mathrm{j}}=\mathrm{m}_{\mathrm{j} 1}+\mathrm{m}_{\mathrm{j} 2}}} c_{\mathrm{m}_{\mathrm{j} 1} \mathrm{~m}_{\mathrm{j} 2}}\left|\mathrm{j}_{1} \mathrm{~m}_{\mathrm{j} 1} ; \mathrm{j}_{2} \mathrm{~m}_{\mathrm{j} 2}\right\rangle, \tag{4.43}
\end{equation*}
$$

where $\mathrm{c}_{\mathrm{mj} 1 \mathrm{~m} 2}$ are Clebsh-Gordan coefficients (or Wigner coefficients). Note, $\mathrm{m}_{\mathrm{j}}=\mathrm{m}_{\mathrm{j} 1}+\mathrm{m}_{\mathrm{j} 2}$. Now, let us determine the coupling constants in case of two spins, as

$$
\begin{equation*}
\left|S M_{S}\right\rangle=\sum_{\substack{m_{s} m_{s} 2 \\ M_{\mathrm{s}}=\mathrm{m}_{\mathrm{s} 1}+\mathrm{m}_{\mathrm{s}}}} c_{\mathrm{m}_{\mathrm{s} 1} \mathrm{~m}_{\mathrm{s} 2}}\left|\mathrm{~m}_{\mathrm{s} 1} \mathrm{~m}_{\mathrm{s} 2}\right\rangle \tag{4.45}
\end{equation*}
$$

Obviously,
$|11\rangle=\alpha_{1} \alpha_{2}$,
thus $\mathrm{c}_{\alpha \alpha}=1$. Now, by using the lowering operator $\mathrm{S}^{-}=\mathrm{s}_{1}^{-}+\mathrm{s}_{2}{ }^{-}$, Eq. (4.30)

$$
\begin{equation*}
\mathrm{S}^{-}\left|\mathrm{S}, \mathrm{M}_{\mathrm{S}}\right\rangle=\hbar \sqrt{\mathrm{S}(\mathrm{~S}+1)-\mathrm{M}_{\mathrm{S}}\left(\mathrm{M}_{\mathrm{S}}-1\right)}\left|\mathrm{IS}, \mathrm{M}_{\mathrm{S}}-1\right\rangle, \tag{4.46}
\end{equation*}
$$

we get from left hand side $S^{-}|11\rangle=\hbar \sqrt{2}|10\rangle$.
From the right we get $\left(\mathrm{s}_{1}^{-}+\mathrm{s}_{2}^{-}\right) \alpha_{1} \alpha_{2}=\mathrm{s}_{1}^{-} \alpha_{1} \alpha_{2}+\mathrm{s}_{2}^{-} \alpha_{1} \alpha_{2}$

$$
\left.=\hbar\left(\beta_{1} \alpha_{2}+\alpha_{1} \beta_{2}\right), \text { thus } \hbar \sqrt{2} 110\right\rangle=\hbar\left(\alpha_{1} \beta_{2}+\beta_{1} \alpha_{2}\right)
$$

which means

$$
\begin{equation*}
|10\rangle=\frac{1}{\sqrt{2}}\left(\alpha_{1} \beta_{2}+\beta_{1} \alpha_{2}\right) . \tag{4.47}
\end{equation*}
$$

Further operation with $\mathrm{S}^{-}=\mathrm{s}_{1}^{-}+\mathrm{s}_{2}^{-}$gives

$$
\begin{equation*}
|1-1\rangle=\beta_{1} \beta_{2} . \tag{4.48}
\end{equation*}
$$

The state 100$\rangle=\mathrm{a} \alpha_{1} \beta_{2}+\mathrm{b} \beta_{1} \alpha_{2}$ is found from the orthogonality condition $\langle 00 \mid 10\rangle=$ ', which together with normalization gives

$$
\begin{equation*}
100\rangle=\frac{1}{\sqrt{2}}\left(\alpha_{1} \beta_{2}-\beta_{1} \alpha_{2}\right) . \tag{4.49}
\end{equation*}
$$

Table of the coupling coefficients for two spins $s_{1}=1 / 2$ and $s_{2}=1 / 2$ :

| $\mathrm{m}_{\mathrm{s} 1}$ | $\mathrm{~m}_{\mathrm{s} 2}$ | $\|11>\|10>\| 0$ | $0>\mid 1-1>$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $\alpha$ | $\alpha$ | 1 | 0 | 0 | 0 |
| $\alpha$ | $\beta$ | 0 | $(1 / 2)^{1 / 2}$ | $(1 / 2)^{1 / 2}$ | 0 |
| $\beta$ | $\alpha$ | 0 | $(1 / 2)^{1 / 2}-(1 / 2)^{1 / 2}$ | 0 |  |
| $\beta$ | $\beta$ | 0 | 0 | 0 | 1 |

### 4.13. Coupling of several angular momenta

Coupling three or more angular momenta can be carried out pairwise, e.g., first two of those, and then, adding a third one, then fourth, and so on.

Example 4.2 Coupling three of the hydrogen atom p-orbitals $\ell_{1,2,3}=1$.

## 5. Group theory

Use of the symmetry properties of the system usually helps in finding the solutions to the Schrödinger equation. It simplifies calculations particularly in cases, where qualitative information is sufficient, e.g., showing if some quantities are equal or zero. Examples of these are the matrix elements (integrals) $\langle\mathrm{n}| \Omega|\mathrm{n}\rangle$ and $\langle\mathrm{n}| \Omega|\mathrm{m}\rangle$ and degeneracy. Also, group theory usually defines the most fundamental and natural (symmetry related) quantum numbers, like those of atomic orbitals: s, p, d, f, g, . ...

Example What are the conditions for $\ell^{2}$ or $\ell_{\mathrm{z}}$ to commute with the Hamiltonian (and therefore, to give good quantum numbers for the energy eigenstates)?
$\mathrm{H}=-\frac{\hbar^{2}}{2 \mathrm{~m}} \nabla^{2}+\mathrm{V}$ and $\ell_{\mathrm{z}}=\frac{\hbar}{\mathrm{i}} \frac{\partial}{\partial \phi}$, thus
$\left[\mathrm{H}, \ell_{\mathrm{z}}\right]=$

## Symmetries of objects

Systematic inspection of symmetry properties is based on the relevant symmetry operations and the group they form. Then the group theory immediately classifies the system and its basic properties.

### 5.1. Symmetry operations and elements

Symmetry operation is an act of doing something to the system, which remains apparently unchanged, e.g., after changing positions of identical parts of the system. The more the system, e.g. a molecule, has symmetry operations the higher is its symmetry. Usual symmetry operatios are rotation, translation, reflection and inversion, which happen with respect to a symmetry element (axis, direction, plane or point).
The symmetry operations of molecules and atoms form sc. point groups (pisteryhmä), which do not include translations as do the more general space groups (avaruusryhmä).
Symmetry operations of point groups always leave one point (center of the molecule) unchanged.

Point groups incorporate five different operations:
E, identity or unit operator (yksikkö- tai ykkösoperaattori), which does nothing.
$\mathrm{C}_{\mathrm{n}}, \mathrm{n}$-fold rotation (n-lukuinen rotaatio) is a rotation by an angle $360^{\circ} / \mathrm{n}$ wrt. symmetry axis. The clockwise (seen from below) rotation is $\mathrm{C}_{\mathrm{n}}{ }^{+}$and counterclockwise $\mathrm{C}_{\mathrm{n}}{ }^{-}$. Note, that $\mathrm{C}_{2}{ }^{+}=\mathrm{C}_{2}{ }^{-}$. If the object has several rotation axes, the one with largest value of n is called principal axis (pääakseli).
$\sigma$, reflection (heijastus) wrt mirror plane (heijastustaso). If the principal axis is in the plane, it is called vertical plane and the operation is denoted by $\sigma_{v}$, and if it is perpendicular to the plane, it is called horizontal plane and the operation is $\sigma_{\mathrm{h}}$. Dihedral plane $\sigma_{\mathrm{d}}$ (and reflection) is a special case of vertical reflection plane, where it bisects the angle between two $\mathrm{C}_{2}$-axes, which are perpendicular to the principal axis.
i, inversion (inversio) wrt center of symmetry. In inversion wrt origin each point of the object at $(\mathrm{x}, \mathrm{y}, \mathrm{z})$ becomes projected to the opposite side with same distance, i.e. to $(-\mathrm{x},-\mathrm{y},-\mathrm{z})$.
$\mathrm{S}_{\underline{n}}$, improper rotation or rotary-reflection (kiertoheijastus) wrt rotation axis. Improper rotation is a composite operation consisting of an $n$-fold rotation and horizontal reflection. Note that $\mathrm{S}_{1}=\sigma_{\mathrm{h}}$ and $\mathrm{S}_{2}=\mathrm{i}$.


Fig. 5.3.



Fig. 7.4. An axis of improper ro-

QTES, 2021

### 5.2. Classification of molecules

The list of all symmetry operations of a molcecule defines its point group. Molecules with the same point group have several common properties. This is an incomplete list of most important symmetry operations:
$\mathrm{C}_{1}$ contains only identity E or 1 .
$\mathrm{C}_{\mathrm{s}}$ : E and one reflection $\sigma$ or m .
$\mathrm{C}_{\mathrm{i}}$ : E and inversion i or 1.
$\mathrm{C}_{\mathrm{n}}$ : E and n -fold rotation $\mathrm{C}_{\mathrm{n}}$ or n .
$\mathrm{C}_{\mathrm{nv}}: \mathrm{E}, \mathrm{C}_{\mathrm{n}}$ and n vertical reflections $\sigma_{\mathrm{v}}$.
$\mathrm{C}_{\mathrm{nh}}: \mathrm{E}, \mathrm{C}_{\mathrm{n}}$ and $\sigma_{\mathrm{h}}$.
$D_{n}: E, C_{n}$ and $n$ two-fold rotations $C_{2}$ perpendicular to $C_{n}$.
$D_{n h}$ : All operations in group $D_{n}$ added by $\sigma_{h}$.
$D_{n d}$ : All operations in group $D_{n}$ and $n$ reflections $\sigma_{d}$.
$\mathrm{S}_{\mathrm{n}}$ : E and $\mathrm{S}_{\mathrm{n}}$. Note, for odd $\mathrm{n} \mathrm{S}_{\mathrm{n}}=\mathrm{C}_{\mathrm{nh}}$.
Note that some combinations of symmetry operations produce new ones, e.g.,
point group $\mathrm{C}_{2 \mathrm{~h}}$ has inversion $\mathrm{i}=\sigma_{\mathrm{h}} \mathrm{C}_{2}$.
Groups T and O are sc. cubic point groups without principal axis. They contain several axes of highest n, similarly as the icosahedral group I.

T: E, $3 \mathrm{C}_{2}, 4 \mathrm{C}_{3}$ and $4 \mathrm{C}_{3}{ }^{\prime}$.
$\mathrm{T}_{\mathrm{d}}$ : "T" $+6 \sigma_{\mathrm{d}}$ and $6 \mathrm{~S}_{4}$ (group of regular tetrahedron).
$T_{h}: T_{d}+i$ (inversion).
$\mathrm{O}: \mathrm{E}, 8 \mathrm{C}_{3}, 3 \mathrm{C}_{2}=3 \mathrm{C}_{4}{ }^{2}$, $6 \mathrm{C}_{2}{ }^{\prime}$ and $6 \mathrm{C}_{4}$.
$\mathrm{O}_{\mathrm{h}}: \mathrm{O}+$ reflections of octahedron (group of regular octahedron).
I: Group of icosahedron.


| $n=$ | 2 | 3 | 4 | 5 | 6 | $\infty$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $C_{n}$ | $1 \underbrace{1}$ |  | $\square$ |  |  |  |
| $D_{n}$ |  |  |  |  |  |  |
| $C_{n v}$ | $\theta$ |  | $\boxed{x}$ |  |  |  |
| $C_{\text {nh }}$ |  | $2$ |  |  |  |  |
| $D_{\text {ah }}$ |  |  |  |  |  |  |
| $D_{n d}$ |  | $\wedge$ | $\Sigma$ |  |  |  |
| $S_{2 n}$ |  | $\lambda$ | $\pi$ |  |  |  |

Fig. 5.18.

QTES, 2021
Above, we have used sc. Schoenflies notation/names for the point groups. There is another notation, which lists the symmetry elements, which is sc. Herman-Mauguin or International notation, in this table:

| $C_{C}: \overline{1}$ | $C_{5}: m$ |  |  |  |
| :--- | :--- | :--- | :--- | :--- |
| $C_{1}: 1$ | $C_{2}: 2$ | $C_{3}: 3$ | $C_{4}: 4$ | $C_{6}: 6$ |
|  | $C_{2 v}: 2 m m$ | $C_{3 v}: 3 m$ | $C_{4 v}: 4 m m$ | $C_{66}: 6 m m$ |
|  | $C_{2 h}: 2 / m$ | $C_{3 h}: \overline{6}$ | $C_{4 h}: 4 / m$ | $C_{6 h}: 6 / m$ |
|  | $D_{2}: 222$ | $D_{3}: 32$ | $D_{4}: 422$ | $D_{6}: 622$ |
|  | $D_{2 h}: m m m$ | $D_{3 h}: \overline{6} 2 m$ | $D_{4 h}: 4 / m m m$ | $D_{6 h}: 6 / m m m$ |
|  | $D_{2 d}: \overline{4} 2 m$ | $D_{3 d}: \overline{3} m$ | $S_{4}: \overline{4}$ | $S_{6}: \overline{3}$ |
| $T: 23$ | $T_{d}: \overline{4} 3 m$ | $T_{h}: m 3$ | $O: 43$ | $O_{h}: m 3 m$ |

The point group of atoms is $R_{3}$, the point group of the sphere, the full spherical symmetry. Naturally, none of the molecules belong to this point group. Properties of $R_{3}$ are the ones of angular momentum.

## Calculus of symmetry

### 5.3. Definition of group

A set of elements with "multiplication"

$$
\mathrm{RS}=\mathrm{T},
$$

where, $\mathrm{R}, \mathrm{S}$ and T belong to this set, form a group, if
(1) identity E belongs to this set,
(2) multiplication is associative, which means that $T(S R)=(T S) R$,
(3) the "product" of two elements also belongs to this set, and
(4) each element $R$ has its "inverse" $R^{-1}$, which also belongs to this set.

Note! The product needs not to be commutative, TS $\neq$ ST .
Note!

$$
\mathrm{R}^{-1}=\mathrm{R}^{-1} \mathrm{R}=\mathrm{E}
$$

It is easy to inspect and find, that the symmetry operations of "point groups" form a groups.

### 5.4. Group multiplication tables

Consider next the representation (esitys) of a point group by matrices, which allows calculation of a "product" $\mathrm{i}=\sigma_{\mathrm{h}} \mathrm{C}_{2}$ by the usual rules of matrix multiplication.
Let us first find multiplication table of one specific group, independent of matrices, as an example.

Consider point group $\mathrm{C}_{3 \mathrm{v}}$ as an example. Now, the symmetry elements of $\mathrm{C}_{3 \mathrm{v}}$ are rotation axis $\mathrm{C}_{3}$ and 3 reflection planes $\sigma_{\mathrm{v}}$, and symmetry operations are $\mathrm{E}, \mathrm{C}_{3}{ }^{+}, \mathrm{C}_{3}{ }^{-}, \sigma_{\mathrm{v}}, \sigma_{\mathrm{v}}{ }^{\prime}$ and $\sigma_{\mathrm{v}}{ }^{\prime \prime}$.
Number of elements in the group is 6 , which is called order of the group, h.
The multiplication table is (shown as $\mathrm{R}=\mathrm{ST}$ )
Table 5.2.

| $\mathrm{S}=\mathrm{T}^{\mathrm{T}}$ | E | $\mathrm{C}_{3}{ }^{+}$ | $\mathrm{C}_{3}{ }^{-}$ | $\sigma_{\mathrm{v}}$ | $\sigma_{\mathrm{v}}{ }^{\prime}$ | $\sigma_{\mathrm{v}}{ }^{\prime \prime}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| E | E | $\mathrm{C}_{3}{ }^{+}$ | $\mathrm{C}_{3}{ }^{-}$ | $\sigma_{\mathrm{v}}$ | $\sigma_{\mathrm{v}}{ }^{\prime}$ | $\sigma_{\mathrm{v}}{ }^{\prime \prime}$ |
| $\mathrm{C}_{3}{ }^{+}$ | $\mathrm{C}_{3}{ }^{+}$ | $\mathrm{C}_{3}{ }^{-}$ | E | $\sigma_{\mathrm{v}}{ }^{\prime}$ | $\sigma_{\mathrm{v}}{ }^{\prime \prime}$ | $\sigma_{\mathrm{v}}$ |
| $\mathrm{C}_{3}{ }^{-}$ | $\mathrm{C}_{3}{ }^{-}$ | E | $\mathrm{C}_{3}{ }^{+}$ | $\sigma_{\mathrm{v}}{ }^{\prime \prime}$ | $\sigma_{\mathrm{v}}$ | $\sigma_{\mathrm{v}}{ }^{\prime}$ |
| $\sigma_{\mathrm{v}}$ | $\sigma_{\mathrm{v}}$ | $\sigma_{\mathrm{v}}{ }^{\prime \prime}$ | $\sigma_{\mathrm{v}}{ }^{\prime}$ | E | $\mathrm{C}_{3}{ }^{-}$ | $\mathrm{C}_{3}{ }^{+}$ |
| $\sigma_{\mathrm{v}}{ }^{\prime}$ | $\sigma_{\mathrm{v}}{ }^{\prime}$ | $\sigma_{\mathrm{v}}$ | $\sigma_{\mathrm{v}}{ }^{\prime \prime}$ | $\mathrm{C}_{3}{ }^{+}$ | E | $\mathrm{C}_{3}{ }^{-}$ |
| $\sigma_{\mathrm{v}}{ }^{\prime \prime}$ | $\sigma_{\mathrm{v}}{ }^{\prime \prime}$ | $\sigma_{\mathrm{v}}{ }^{\prime}$ | $\sigma_{\mathrm{v}}$ | $\mathrm{C}_{3}{ }^{-}$ | $\mathrm{C}_{3}{ }^{+}$ | E |

Fig. 5.19.

### 5.5. Matrix representations

Matrices, which obey this multiplication table can be chosen in many $(\infty)$ different ways. The choice is defined by sc. basis, which we now choose as shown in the figure. The picture shows, e.g. 1 s -orbitals of the atoms in $\mathrm{NH}_{3}$ molecule.

Fig. 5.23.

Dimension of this basis is 4, i.e. the number of basis functions. The basis can be given as a "vector" $\mathbb{f}=\left(\mathrm{s}_{\mathrm{N}}, \mathrm{s}_{\mathrm{A}}, \mathrm{s}_{\mathrm{B}}, \mathrm{s}_{\mathrm{C}}\right)$, and the effect of symmetry operations can be written as $\sigma_{\mathrm{v}}\left(\mathrm{s}_{\mathrm{N}}, \mathrm{s}_{\mathrm{A}}, \mathrm{s}_{\mathrm{B}}, \mathrm{s}_{\mathrm{C}}\right)=\left(\mathrm{s}_{\mathrm{N}}, \mathrm{s}_{\mathrm{A}}, \mathrm{s}_{\mathrm{C}}, \mathrm{s}_{\mathrm{B}}\right)$. This can be written as a matrix multiplication

$$
=\left(\mathrm{s}_{\mathrm{N}}, \mathrm{~s}_{\mathrm{A}}, \mathrm{~s}_{\mathrm{B}}, \mathrm{~s}_{\mathrm{C}}\right)\left[\begin{array}{llll} 
& \sigma_{\mathrm{v}}\left(\mathrm{~s}_{\mathrm{N}}, \mathrm{~s}_{\mathrm{A}}, \mathrm{~s}_{\mathrm{B}}, \mathrm{~s}_{\mathrm{C}}\right)=  \tag{5.3}\\
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 0 & 1 \\
0 & 0 & 1 & 0
\end{array}\right]=\left(\mathrm{s}_{\mathrm{N}}, \mathrm{~s}_{\mathrm{A}}, \mathrm{~s}_{\mathrm{C}}, \mathrm{~s}_{\mathrm{B}}\right) .
$$

This matrix is called a representation (esitys) of $\sigma_{v}$ and it is denoted by $\mathbb{D}\left(\sigma_{v}\right)$, whose components can be written as $\mathrm{D}_{\mathrm{ji}}\left(\sigma_{\mathrm{v}}\right)$. One should note, that multiplication by a matrix only means that

$$
\begin{array}{ll}
\sigma_{v} s_{N}=s_{N}, & \sigma_{v} s_{A}=s_{A}, \\
\sigma_{v} s_{B}=s_{C} j a & \sigma_{v} s_{C}=s_{B} .
\end{array}
$$

Generally for any basis function $f_{i}$ and operation $R$, we can write

$$
\begin{equation*}
\mathrm{Rf} \mathrm{f}_{\mathrm{i}}=\sum_{\mathrm{j}} \mathrm{f}_{\mathrm{j}} \mathrm{D}_{\mathrm{ji}}(\mathrm{R}) . \tag{5.4}
\end{equation*}
$$

Similarly for $\mathrm{C}_{3}{ }^{+}$

$$
\mathrm{C}_{3}^{+}\left(\mathrm{s}_{\mathrm{N}}, \mathrm{~s}_{\mathrm{A}}, \mathrm{~s}_{\mathrm{B}}, \mathrm{~s}_{\mathrm{C}}\right)=
$$

$$
=\left(\mathrm{s}_{\mathrm{N}}, \mathrm{~s}_{\mathrm{A}}, \mathrm{~s}_{\mathrm{B}}, \mathrm{~s}_{\mathrm{C}}\right)\left[\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0
\end{array}\right]=\left(\mathrm{s}_{\mathrm{N}}, \mathrm{~s}_{\mathrm{B}}, \mathrm{~s}_{\mathrm{C}}, \mathrm{~s}_{\mathrm{A}}\right) \quad \text { or } \quad \mathrm{C}_{3}^{+} \mathrm{f}_{\mathrm{i}}=\sum_{\mathrm{j}} \mathrm{f}_{\mathrm{j}} \mathrm{D}_{\mathrm{ji}}\left(\mathrm{C}_{3}^{+}\right) \text {, }
$$

for any $\mathrm{f}_{\mathrm{i}}$. The representations of all operations in group $\mathrm{C}_{3 \mathrm{v}}$ are given in basis ( $\mathrm{s}_{\mathrm{N}}, \mathrm{s}_{\mathrm{A}}, \mathrm{s}_{\mathrm{B}}$, $\mathrm{s}_{\mathrm{C}}$ ) in the Table 5.3, on next page.
Consider now two consequent operations or "product", e.g. $\sigma_{v} \mathrm{C}_{3}{ }^{+}=\sigma_{v} "$ with matrices.

QTES, 2021

$$
\begin{align*}
& \begin{array}{l}
\mathrm{D}\left(\sigma_{\mathrm{v}}\right) \mathrm{D}\left(\mathrm{C}_{3}^{+}\right)= \\
=\left[\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0 \\
0 & 0 & 0 & 1 \\
0 & 0 & 1 & 0
\end{array}\right]\left[\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 \\
0 & 1 & 0 & 0 \\
0 & 0 & 1 & 0
\end{array}\right]=\left[\begin{array}{llll}
1 & 0 & 0 & 0 \\
0 & 0 & 0 & 1 \\
0 & 0 & 1 & 0 \\
0 & 1 & 0 & 0
\end{array}\right]
\end{array} \\
& =\mathrm{D}\binom{\mathrm{o}}{\mathrm{v}} \\
& \mathbb{D}(\mathrm{R}) \mathbb{D}(\mathrm{S})=\mathbb{D}(\mathrm{RS}) \text {. } \tag{5.5}
\end{align*}
$$

This means that the "product" of operations can be found by matrix multiplication. In general, if $R$ and $S$ are symmetry operations of a group, then

Table 5.3. Matrix representations of $C_{3 v}$ in the basis $\left(s_{N}, s_{A}, s_{B}, s_{C}\right)$.

| D(E) | $\mathrm{D}\left(\mathrm{C}_{3}^{+}\right)$ | $\mathrm{D}\left(\mathrm{C}_{3}^{-}\right)$ |  |
| :---: | :---: | :---: | :---: |
| $\begin{array}{cccc}1 & 0 & 0 & 0\end{array}$ | $\begin{array}{lllll}1 & 0 & 0 & 0\end{array}$ | $\begin{array}{lllll}1 & 0 & 0 & 0\end{array}$ | of two groups is the |
| $\begin{array}{lllll}0 & 1 & 0 & 0\end{array}$ | $\begin{array}{lllll}0 & 0 & 0 & 1\end{array}$ | $\begin{array}{llll}0 & 0 & 1 & 0\end{array}$ | same, the groups are |
| $\begin{array}{lllll}0 & 0 & 1 & 0\end{array}$ | $\begin{array}{lllll}0 & 1 & 0 & 0\end{array}$ | $\begin{array}{lllll}0 & 0 & 0 & 1\end{array}$ | called homomorphic. |
| $\begin{array}{lllll}0 & 0 & 0 & 1\end{array}$ | $\begin{array}{llll}0 & 0 & 1 & 0\end{array}$ | $\begin{array}{llll}0 & 1 & 0 & 0\end{array}$ |  |
| $\chi(\mathrm{E})=4$ | $\chi\left(\mathrm{C}_{3}^{+}\right)=1$ <br> ja | $\chi\left(\mathrm{C}_{3}^{-}\right)=1$ | Note, that notations are $(\mathrm{RS}) \mathbb{f}=\mathfrak{f} \mathbb{D}(\mathrm{RS})$ |
| $\mathrm{D}\left(\sigma_{\mathrm{v}}\right)$ | $\mathrm{D}\left(\sigma_{\mathrm{v}}^{\prime}\right)$ | $\mathrm{D}\left(\sigma_{\mathrm{v}}^{\prime \prime}\right)$ | $=\mathfrak{f}[\mathbb{D}(\mathrm{R}) \mathbb{D}(\mathrm{S})]$ |
| $\begin{array}{lllll}1 & 0 & 0 & 0\end{array}$ | $\begin{array}{llll}1 & 0 & 0 & 0\end{array}$ | $\begin{array}{llll}1 & 0 & 0 & 0\end{array}$ | $[f \mathbb{D}(\mathrm{R})] \mathbb{D}(\mathrm{S})$, |
| $\begin{array}{lllll}0 & 1 & 0 & 0\end{array}$ | $\begin{array}{lllll}0 & 0 & 1 & 0\end{array}$ | $\begin{array}{lllll}0 & 0 & 0 & 1\end{array}$ | though |
| $\begin{array}{lllll}0 & 0 & 0 & 1\end{array}$ | $\begin{array}{lllll}0 & 1 & 0 & 0\end{array}$ | $\begin{array}{llll}0 & 0 & 1 & 0\end{array}$ | (RS) $\mathfrak{f f}=\mathrm{R}(\mathrm{Sff}$; |
| $\left[\begin{array}{lllll}0 & 0 & 1 & 0\end{array}\right]$ | $\begin{array}{llll}0 & 0 & 0 & 1\end{array}$ | $\left[\begin{array}{llll}0 & 1 & 0 & 0\end{array}\right.$ |  |
| $\chi\left(\sigma_{v}\right)=2$ | $\chi\left(\sigma_{\mathrm{v}}^{\prime}\right)=2$ | $\chi\left(\sigma_{v}^{\prime \prime}\right)=2$ |  |

### 5.6. Properties of matrix representations

Let us consider another basis $\mathbb{f}^{\prime}=\left(\mathrm{s}_{\mathrm{N}}, \mathrm{s}_{1}, \mathrm{~s}_{2}, \mathrm{~s}_{3}\right)$,
which is defined by $f=\left(s_{N}, s_{A}, s_{B}, s_{C}\right)$ such, that
$\mathrm{s}_{1}=\mathrm{s}_{\mathrm{A}}+\mathrm{s}_{\mathrm{B}}+\mathrm{s}_{\mathrm{C}}, \quad \mathrm{s}_{2}=2 \mathrm{~s}_{\mathrm{A}}-\mathrm{s}_{\mathrm{B}}-\mathrm{s}_{\mathrm{C}}$ and $\mathrm{s}_{3}=\mathrm{s}_{\mathrm{B}}-\mathrm{s}_{\mathrm{C}}$.
The transformation can be written as

$$
\mathrm{f}_{\mathrm{i}}^{\prime}=\sum_{\mathrm{j}} \mathrm{f}_{\mathrm{j}} \mathrm{c}_{\mathrm{ji}}
$$

or

$$
\begin{equation*}
\mathfrak{f}^{\prime}=\mathfrak{f} \mathfrak{c}, \tag{*}
\end{equation*}
$$

where

$$
\mathbb{C}=\left[\begin{array}{cccc}
1 & 0 & 0 & 0 \\
0 & 1 & 2 & 0 \\
0 & 1 & -1 & 1 \\
0 & 1 & -1 & -1
\end{array}\right]
$$

Equation (5.4) can be written in compact form

$$
\begin{equation*}
\mathrm{R} \mathbb{f}=\mathbb{f} \mathbb{D}(\mathrm{R}) \tag{**}
\end{equation*}
$$

and correspondingly, in basis $\mathfrak{f}^{\prime}$,

$$
R \mathbb{f}^{\prime}=f^{\prime} \mathbb{D}^{\prime}(\mathrm{R}) .
$$

With substitution of (*), we get $R \mathbb{f} c=f \in \mathbb{D}(\mathrm{R})$ and with multiplication from right by $c^{-1}$, follows $\mathrm{R} f=f \in \mathbb{D}^{\prime}(\mathrm{R}) \mathbb{c}^{-1}$. Comparison with $\left({ }^{* *}\right)$ allows us to write sc. similarity transformation

$$
\begin{equation*}
\mathbb{D}(R)=\mathbb{C} \mathbb{D}^{\prime}(\mathrm{R}) \mathbb{c}^{-1} \tag{5.7a}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbb{D}^{\prime}(\mathrm{R})=\mathbb{c}^{-1} \mathbb{D}(\mathrm{R}) \mathbb{c} . \tag{5.7b}
\end{equation*}
$$

QTES, 2021
As

$$
c^{-1}=\left[\begin{array}{cccc}
6 & 0 & 0 & 0 \\
0 & 2 & 2 & 2 \\
0 & 2 & -1 & -1 \\
0 & 0 & 3 & -3
\end{array}\right] / 6,
$$

the matrix representation of $\mathrm{C}_{3 \mathrm{v}}$ in basis $f^{\prime}$ can be calculated.
Table 5.4. Matrix representation of $\mathrm{C}_{3 \mathrm{v}}$ in basis ( $\mathrm{s}_{\mathrm{N}}, \mathrm{s}_{1}, \mathrm{~s}_{2}, \mathrm{~s}_{3}$ ).


In Tables 5.3 and 5.4 it is also given the trace (jälki, German Spur) of the matrices, which is the sum of diagonal elements. Compare!

### 5.7. Characters of representations

The trace of the matrix (sum of the diagonal elements) is its character (karakteeri)

$$
\begin{equation*}
\chi(\mathrm{R})=\sum_{\mathrm{i}} \mathrm{D}_{\mathrm{ii}}(\mathrm{R})=\operatorname{Tr} \mathrm{D}(\mathrm{R}) \tag{5.8}
\end{equation*}
$$

We note, that the characters of symmetry operations remain in similarity transformations. This is easy to generalise, because

$$
\begin{equation*}
\operatorname{Tr} \mathbb{A B C}=\operatorname{Tr} \mathbb{B C A}=\operatorname{Tr} \mathbb{C} \mathbb{A} B \tag{5.9}
\end{equation*}
$$

Thus, in general

$$
\chi(\mathrm{R})=\chi^{\prime}(\mathrm{R}) .
$$

### 5.8. Characters and classes

In addition we note, that different types of operations have different and similar ones have the same characters. Now, for $\mathrm{C}_{3 \mathrm{v}}$ we have: $\chi(\mathrm{E})=4, \chi($ rotation $)=1$ and $\chi$ (reflection $)=2$. Thus, we see that the character characterizes or classifies the elements of the group, and therefore, we define the concept class (luokka):
The elements $\mathrm{R}_{1}$ and $\mathrm{R}_{2}$ belong to the same class, if such element S belongs to the group, that

$$
\begin{equation*}
\mathrm{R}_{1}=\mathrm{S}^{-1} \mathrm{R}_{2} \mathrm{~S} . \tag{5.10}
\end{equation*}
$$

However, depending on the representation it may accidentally happen that operations belonging to different classes have the same character. As an example, consider the one dimensional representation consisting of $1 \times 1$ matrices [1], only.

### 5.9. Irreducible representations

If the matrix representation of the goup is or can be transformed to block-diagonal form, like the matrices in Table 5.3, the representation can be reduced to two (or more) representations. Thus, the 4-dimensional representation in Table 5.3 can be reduced to 1- and 3dimensional representations, as $\mathbb{D}^{(4)}=\mathbb{D}^{(1)} \oplus \mathbb{D}^{(3)}$, (direct sum). The 1-dimensional representation consists of six $1 \times 1$ matrices [1], which obey the multiplication table of $\mathrm{C}_{3 \mathrm{v}}$, Table 5.2.
After transformation to another basis the reduction can be continued, as seen in Table 5.4. We obtain $\mathbb{D}^{(4)}=2 \mathbb{D}^{(1)} \oplus \mathbb{D}^{(2)}$. The reduction can not be continued further, and therefore, these are the irreducible representations of the point group $\mathrm{C}_{3 \mathrm{v}}$ (redusoitumaton esitys). From now on let us use a short form "irrep".

Let us make a look at the basis sets of irreps. We see that the functions $\mathrm{s}_{\mathrm{N}}$ and $\mathrm{s}_{1}$ have the same symmetry species (symmetrialaji) in group $\mathrm{C}_{3 \mathrm{v}}$ and they span (virittää) or form the bases for the two 1-dimensional irreps. Functions $\mathrm{s}_{2}$ and $\mathrm{s}_{3}$ are of different symmetry (species) and span the 2-dim irrep.

Thus, we learn that "different" basis functions span different irreps, and therefore, irreps can be used to "describe symmetry properties".

Let us name the irreps or symmetry species. The 1-dim irrep with characters ( $1,1,1,1,1,1$ ) is called $\Gamma_{1}$ and 2-dim irrep with characters $(2,-1,-1,0,0,0)$ is called $\Gamma_{3}$. The more usual notations for these are $\mathrm{A}_{1}$ and E . Note that incidentally the identity operation is denoted by E , too.

Figs. 5.24.

### 5.10. Orthogonality theorems

For irreducible representations we have the great orthogonality theorem (GOT):

$$
\begin{equation*}
\sum_{\mathrm{R}} \mathrm{D}_{\mathrm{ij}}^{f}(\mathrm{R})^{*} \mathrm{D}_{\mathrm{ij}} \epsilon^{\prime}(\mathrm{R})=\frac{\mathrm{h}}{\mathrm{~d}_{\ell}} \delta_{\ell \ell^{\prime}} \delta_{\mathrm{ii}} \delta_{\mathrm{jj}}, \tag{5.12}
\end{equation*}
$$

where h is the order of the group, $\ell$ refers to the irrep $\Gamma_{\ell}$ and $\mathrm{d}_{\ell}$ is the dimension of the irrep. For irreps we also have the little orthogonality theorem (LOT):

$$
\begin{equation*}
\sum_{\mathrm{R}} \chi^{\ell}(\mathrm{R})^{*} \chi^{\ell}(\mathrm{R})=\mathrm{h} \delta_{\ell \ell} \text {, or } \quad \sum_{\mathrm{c}} \mathrm{~g}(\mathrm{c}) \chi^{\ell}(\mathrm{c})^{*} \chi^{\ell}(\mathrm{c})=\mathrm{h} \delta_{\ell \ell}, \tag{5.13,5.14}
\end{equation*}
$$

where c refers to the classes and $\mathrm{g}(\mathrm{c})$ is the number of symmetry operations in class c .
The GOT implies that
number of symmetry species $=$ number classes
and

$$
\begin{equation*}
\sum_{l} \mathrm{~d}_{\ell}^{2}=\mathrm{h} . \tag{5.18}
\end{equation*}
$$

Consider $\mathrm{C}_{3 \mathrm{v}}$, which consists 3 classes. Then we have 3 symmetry species (and irreps), of which we know two, $\Gamma_{1}$ and $\Gamma_{3}$ or $\mathrm{A}_{1}$ and E , already. Let us name the third one as $\Gamma_{2}$ or $A_{2}$. Now, with (5.18) we can write $1^{2}+\mathrm{d}_{2}^{2}+2^{2}=6$, which gives $\mathrm{d}_{2}=1$. By using the LOT (5.14) we can now fill the missing row in the character table of $\mathrm{C}_{3 \mathrm{v}}$, Table 5.5.
Note! Also the columns are orthogonal to each other.

Table 5.5.
Character table of group $\mathrm{C}_{3 \mathrm{v}}$.

| $\mathrm{C}_{3 \mathrm{v}}$, | E | $2 \mathrm{C}_{3}$ | $3 \sigma_{\mathrm{v}}$ |
| :---: | :---: | :---: | :---: |
| $\mathrm{A}_{1}$ | 1 | 1 | 1 |
| $\mathrm{~A}_{2}$ |  |  |  |
| E | 2 | -1 | 0 |

QTES, 2021

## Reduced representations

Let us find out how to identify the symmetry species, which are spanned by a given basis.

### 5.11.Reduction of representations

To identify the irreps, which the basis functions span, we need to reduce the matrix rep

$$
\begin{equation*}
\mathbb{D}(\mathrm{R})=\sum_{\ell}{ }^{\oplus} \mathrm{a}_{\ell} \mathbb{D}^{(\ell)}(\mathrm{R}), \tag{5.19}
\end{equation*}
$$

corresponding the symmetry elements $\Gamma_{\ell}$

$$
\begin{equation*}
\Gamma=\sum_{\ell} \mathrm{a}_{\ell} \Gamma_{\ell} . \tag{5.20}
\end{equation*}
$$

E.g., for the basis ( $\mathrm{s}_{\mathrm{N}}, \mathrm{s}_{\mathrm{A}}, \mathrm{s}_{\mathrm{B}}, \mathrm{s}_{\mathrm{C}}$ ) and $\mathrm{C}_{3 \mathrm{v}}$ group we found $\Gamma=2 \mathrm{~A}_{1}+\mathrm{E}$.

Thus, in general, we need to find the "reduction coefficients" $\mathrm{a}_{\ell}$.
Similarity transformation preserves the trace of matrices, and therefore,

$$
\begin{equation*}
\chi(\mathrm{R})=\sum_{\ell} \mathrm{a}_{\ell} \chi^{(\ell)}(\mathrm{R}), \tag{5.21}
\end{equation*}
$$

which together with the LOT implies

$$
\begin{equation*}
\mathrm{a}_{\ell}=\frac{1}{\mathrm{~h}} \sum_{\mathrm{R}} \chi^{(\vartheta)}(\mathrm{R})^{*} \chi(\mathrm{R}) \tag{5.22}
\end{equation*}
$$

and/or

$$
\begin{equation*}
\mathrm{a}_{\ell}=\frac{1}{\mathrm{~h}} \sum_{\mathrm{c}} \mathrm{~g}(\mathrm{c}) \chi^{(\ell)}(\mathrm{c})^{*} \chi(\mathrm{c}) . \tag{5.23}
\end{equation*}
$$

Often, for finding the symmetry species it is sufficient to inspect the character tables, only.

Example 5.8. What symmetry species do the s-orbitals of atoms in $\mathrm{CH}_{4}$ molecule span?

|  | $E$ | $8 C_{3}$ | $3 C_{2}$ | $6 S_{4}$ | $6 \sigma_{\mathrm{d}}$ |
| :--- | :--- | ---: | ---: | ---: | ---: |
| $\mathrm{A}_{1}$ | 1 | 1 | 1 | 1 | 1 |
| $\mathrm{~A}_{2}$ | 1 | 1 | 1 | -1 | -1 |
| E | 2 | -1 | 2 | 0 | 0 |
| $\mathrm{~T}_{1}$ | 3 | 0 | -1 | 1 | -1 |
| $\mathrm{~T}_{2}$ | 3 | 0 | -1 | -1 | 1 |

### 5.12. Symmetry-adapted bases

Next, a method is described for projecting out from a given basis $\mathbb{f}$ a new one sc. symmetryadapted basis $\mathbb{f}^{\prime}$, which spans the irreps of given symmetry species.

A projection operator

$$
\begin{equation*}
P_{i j}^{(i)}=\frac{d_{\epsilon}}{h} \sum_{\mathrm{R}} D_{i j}^{(\hat{l}(\mathrm{R}) * R} \tag{5.24}
\end{equation*}
$$

has the property

$$
\begin{equation*}
\mathrm{P}_{\mathrm{ij}}^{(\hat{i}} \mathrm{f}_{\mathrm{j}}^{(\ell)}=\mathrm{f}_{\mathrm{i}}^{(\ell)} \delta_{\ell \ell} \delta_{\mathrm{j} \mathrm{j}^{\prime}}, \tag{5.25}
\end{equation*}
$$

thus, it projects from the irrep $\ell$ function $\mathrm{f}_{\mathrm{j}}{ }^{(\ell)}$ another function $\mathrm{f}_{\mathrm{i}}{ }^{(\ell)}$ belonging to the same irrep. A projection operator

$$
\begin{equation*}
\mathrm{p}^{(\vartheta)}=\sum_{\mathrm{i}} \mathrm{P}_{\mathrm{ii}}^{(\vartheta)}=\frac{\mathrm{d}_{\ell}}{\mathrm{h}} \sum_{\mathrm{R}} \chi^{(\theta)}(\mathrm{R})^{*} \mathrm{R} \tag{5.29-30}
\end{equation*}
$$

has the property

$$
\begin{equation*}
\mathrm{p}^{(\ell)} \mathrm{f}_{\mathrm{j}}=\sum_{\mathrm{i}} \mathrm{f}_{\mathrm{i}}^{()}, \tag{5.31}
\end{equation*}
$$

in words, it projects from function $\mathrm{f}_{\mathrm{j}}$ a sum of all functions belonging to $\ell$.

Example 5.9. Construct the symmetry-adapted bases from basis ( $\mathrm{s}_{\mathrm{A}}, \mathrm{s}_{\mathrm{B}}, \mathrm{s}_{\mathrm{C}}$ ) for group $\mathrm{C}_{3 \mathrm{v}}$.

## Symmetry properties of functions

### 5.13. Transformation of p-orbitals

Consider the p-orbitals of N in $\mathrm{NH}_{3}$ molecule as a basis for symmetry species in group $\mathrm{C}_{3 \mathrm{v}}$. The real p-orbitals are
$p_{x}=r \sin \theta \cos \phi f(r)=x f(r)$,
$p_{y}=r \sin \theta \sin \phi f(r)=y f(r)$ and
$p_{z}=r \cos \theta f(r)=z f(r)$,
where $f(r)$ is the spherically symmetric radial part. Thus, the symmetry properties are the same as those of the basis $f=(x, y, z)$, whose transformations for the operations in $\mathrm{C}_{3 \mathrm{v}}$ are shown in Fig. 5.28. Thus, we can write, e.g.

$$
\sigma_{\mathrm{v}}(\mathrm{x}, \mathrm{y}, \mathrm{z})=(-\mathrm{x}, \mathrm{y}, \mathrm{z})=(\mathrm{x}, \mathrm{y}, \mathrm{z})\left[\begin{array}{ccc}
-1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 1
\end{array}\right]
$$

and

$$
\mathrm{C}_{3}^{+}(\mathrm{x}, \mathrm{y}, \mathrm{z})=\left(-\frac{1}{2} \mathrm{x}+\frac{1}{2} \sqrt{3} \mathrm{y},-\frac{1}{2} \sqrt{3} \mathrm{x}-\frac{1}{2} \mathrm{y}, \mathrm{z}\right)=(\mathrm{x}, \mathrm{y}, \mathrm{z})\left[\begin{array}{ccc}
-\frac{1}{2} & -\frac{1}{2} \sqrt{3} & 0 \\
\frac{1}{2} \sqrt{3} & -\frac{1}{2} & 0 \\
0 & 0 & 1
\end{array}\right] .
$$


etc. Representation of group $\mathrm{C}_{3 \mathrm{v}}$ in basis $(\mathrm{x}, \mathrm{y}, \mathrm{z})$ is given in Table 5.7.
Because the matrices are in block-diagonal form, we see, that z spans the species $\mathrm{A}_{1}$ and functions (x,y) span the species E. Note, that the sum of characters add up correctly

$$
(3,0,1)=(1,1,1)+(2,-1,0) .
$$

Fig. 5.29.


### 5.14. Direct-product bases and atomic d-orbitals

Table 5.7

A product of two bases with dimensions $\mathrm{d}_{1}$ and $\mathrm{d}_{2}$ form a sc. direct-product basis, whose dimension is $d=d_{1} d_{2}$. For a direct product of two bases, which span $\Gamma_{\ell}$ and $\Gamma_{\ell^{\prime}}$, it can be shown that

$$
\begin{equation*}
\chi(\mathrm{R})=\chi^{(\ell)}(\mathrm{R}) \chi^{\left(\ell^{\prime}\right)}(\mathrm{R}) . \tag{5.32}
\end{equation*}
$$

Direct-product of basis ( $\mathrm{x}, \mathrm{y}, \mathrm{z}$ ) with itself is $(x, y, z) \times(x, y, z)=$
 ( $\mathrm{x}^{2}, \mathrm{xy}, \mathrm{xz}, \mathrm{yx}, \mathrm{y}^{2}, \mathrm{yz}, \mathrm{zx}, \mathrm{zy}, \mathrm{z}^{2}$ ). Basis ( $\mathrm{x}, \mathrm{y}, \mathrm{z}$ ) spans species $\mathrm{A}_{1}$ and E , and its characters are 3,0 and 1 , see Table 5.7. Now, eq. (5.32) implies characters 9,0 and 1 , and by inspecting the character table of $\mathrm{C}_{3 \mathrm{v}}$ we find, that these numbers can be found from the direct-sum $2 A_{1}+A_{2}+3 \mathrm{E}$, only. Further inspection leads to conclusions, that

$$
\begin{array}{lll}
(\mathrm{z}) \times(\mathrm{x}, \mathrm{y})=(\mathrm{xz}, \mathrm{yz}) & -> & \mathrm{A}_{1} \times \mathrm{E}=\mathrm{E} \\
(\mathrm{x}, \mathrm{y}) \times(\mathrm{x}, \mathrm{y})=\left(\mathrm{x}^{2}, \mathrm{xy}, \mathrm{yx}, \mathrm{y}^{2}\right) & -> & \mathrm{E} \times \mathrm{E}=\mathrm{A}_{1}+\mathrm{A}_{2}+\mathrm{E} .
\end{array}
$$

Direct-products of symmetry species are tabulated, Appendix 1. Further analysis shows, that the symmetry-adapted $x^{2}+y^{2}$ spans $A_{1},\left(x^{2}-y^{2}, x y+y x=2 x y\right)$ span $E$ and ( $x y-y x=0$ ) corresponds $\mathrm{A}_{2}$. Also these can be found in tables in Appendix 1.
This is a direct (and easy) way to find atomic d-orbitals ( $\mathrm{xy}, \mathrm{yz}, \mathrm{zx}, \mathrm{x}^{2}-\mathrm{y}^{2}, 3 \mathrm{z}^{2}-\mathrm{r}^{2}$ ) and the symmetry species they span.

QTES, 2021

### 5.15. Direct-product groups

Consider two groups G and $\mathrm{G}^{\prime}$ of the order h and $\mathrm{h}^{\prime}$, whose only common element is E and whose elements $R_{i}$ and $R_{j}$ ' commute, for $i=1,2, \ldots, h$ and $j=1,2, \ldots, h^{\prime}$. Symmetry operations $R_{i} R_{j}^{\prime}$ form the direct-product group of groups $G$ and $G^{\prime}$ and it is denoted by $\mathrm{G}^{\prime \prime}=\mathrm{G} \otimes \mathrm{G}^{\prime}$. For the characters it holds that

$$
\begin{equation*}
\chi\left(\mathrm{RR}^{\prime}\right)=\chi(\mathrm{R}) \chi\left(\mathrm{R}^{\prime}\right), \tag{5.37}
\end{equation*}
$$

h" = hh' and similarly for the number of classes: \#(classes)" = \#(classes) x \#(classes)'.

Example 5.11. Construct the character table of direct-product group $D_{3 h}=C_{3 v} \otimes C_{s}$ using (5.37)

| $C_{3 v}$ <br> $(3 m)$ | $E$ | $2 C_{3}$ | $3 \sigma_{v}$ |
| :--- | :---: | ---: | ---: |
| $\mathrm{~A}_{1}$ | 1 | 1 | 1 |
| $\mathrm{~A}_{2}$ | 1 | 1 | -1 |
| E | 2 | -1 | 0 |
|  |  |  |  |
| $C_{\mathrm{s}}=C_{\mathrm{h}}$ | $E$ | $\sigma_{\mathrm{h}}$ |  |
| $(m)$ |  |  |  |
| $\mathrm{A}^{\prime}$ | 1 | 1 |  |
| $\mathrm{~A}^{\prime \prime}$ | 1 | -1 |  |

### 5.16. Vanishing integrals

Let us consider evaluation of an integral over a range for an odd function $\mathrm{f}: \mathrm{f}(-\mathrm{x})=-\mathrm{f}(\mathrm{x})$ and an even function $\mathrm{g}: \mathrm{g}(-\mathrm{x})=\mathrm{g}(\mathrm{x})$ in range $-\mathrm{a}<\mathrm{x}<\mathrm{a}$. It is easy to see that the integral for the odd function vanishes identically, whereas for the even function it does not, necessarily. However, incidentally the latter one can vanish, as well.

Symmetry analysis of the range ( -a , a) leads to the same result, as follows. The symmetry operations are E and $\sigma_{\mathrm{h}}$, and thus, it belongs to the group $\mathrm{C}_{\mathrm{s}}$. Functions g and f span symmetry species $\mathrm{A}^{\prime}$ and $\mathrm{A}^{\prime \prime}$, in this order. Now, based on the possible symmetry of the integrand, we can conclude, that
an integral of an integrand with certain symmetry over a symmetric range vanishes, except if the symmetry species of the integrand is highest possible, sc. full symmetry, usually $\mathrm{A}_{1}$.
This result can be generalized and used in evaluation of matrix elements, e.g., for expectation values and transition probabilities, etc. It is often called as "matrix element rule".

Fig. 5.32.

QTES, 2021
In case the integrand is a product of several functions and/or operators, its symmetry species is found as a direct-product of its factors. If it is other than the highest (usually $\mathrm{A}_{1}$ ), the integral vanishes, otherwise not necessarily.
An alternative way of finding the same conclusion can be carried out based on the orthogonality analysis

$$
\begin{equation*}
\int_{\tau} \mathrm{f}_{\mathrm{i}}^{(\ell) * f_{\mathrm{j}}^{(\ell)} \mathrm{d} \tau \propto \delta_{\ell \ell} \delta_{\mathrm{ij}} . . . . . . . .} \tag{5.38}
\end{equation*}
$$

Example 5.12. Find those orbitals of N atom in $\mathrm{NH}_{3}$ molecule, which may have nonzero overlap integral (peittointegraali) with the 1 s orbitas of H , i.e., $\mathrm{s}_{1}, \mathrm{~s}_{2}$ and $\mathrm{s}_{3}$.

### 5.17. Symmetry and degeneracy

Hamiltonian of a system has to be invariant (remain unchanged) in all symmetry operations. Therefore, the hamiltonian has the full symmetry and it commutes with all other symmetry operations. This implies, that $\psi$ and $\mathrm{R} \psi$ are both eigenfunctions of the Schrödinger equation $\mathrm{H} \psi=\mathrm{E} \psi$ corresponding to the same eigenenergy E . Thus, the eigenfunctions of the same symmetry species are degenerate with the same eigenenergy and all eigenfunctions of that species can be found by using symmetry operators (or projection operators).
Therefore, the degree of degeneration is the dimension of irrep, $\chi(\mathrm{E})$.

## Full rotation group

### 5.18. Generators of rotations

The point group of heteronuclear diatomic or linear molecule is $\mathrm{C}_{\infty \mathrm{v}}$ and that of homonuclear diatomic is $\mathrm{D}_{\infty h}$. Properties of $\mathrm{C}_{\infty v}$ allows derivation of the properties of the z-komponent of angular moment, i.e., operator $\ell_{z}$. This follows from the fact that rotation by any angle $\phi$ with respect to a fixed axis is one of the symmetry operations of the group.


Fig. 5.35.

QTES, 2021

### 5.19. Point group of sphere

The point group of spherical atoms is $\mathrm{R}_{3}$ and it gives the properties of angular momentum $\ell^{2}$ including the commutation rules. Coupling of two angular momenta can be done by reducing the direct-product of species spanned by the irreps of those two.
Writing the representation of rotation $\mathrm{C}_{\mathrm{a}}{ }^{(\mathrm{z})}$ in a basis of spherical harmonics $\left\{\mathrm{Y}_{\ell^{\prime} \ell}, \mathrm{Y}_{\ell^{\prime} \ell-1}, \ldots, \mathrm{Y}_{\left.\ell^{\prime}-\ell\right\}}\right.$ it can be seen that its character is

$$
\begin{equation*}
\chi\left(\mathrm{C}_{\alpha}\right)=\frac{\sin \left[\left(\ell+\frac{1}{2}\right) \alpha\right]}{\sin \left[\frac{1}{2} \alpha\right]} . \tag{5.47b}
\end{equation*}
$$

## 6. Perturbation theory and variation theorem

In practical calculations, after possibly utilizing group theory, approximate methods are needed for finding the numerical values. Next, we will consider the two most applied ones: perturbation theory and variation theorem, and later, also the iteration algorithm.

## Semiclassical WKB approximation

Let us start from the total energy in classical mechanics, $E=p^{2} / 2 m+V$, and one-dimensional time-independent S-equation

$$
\begin{equation*}
\hbar^{2} \frac{\mathrm{~d}^{2} \psi}{\mathrm{~d} x^{2}}+p^{2} \psi=0 \tag{6.1a}
\end{equation*}
$$

where now

$$
\begin{equation*}
\mathrm{p}(\mathrm{x})=\{2 \mathrm{~m}[\mathrm{E}-\mathrm{V}(\mathrm{x})]\}^{1 / 2} . \tag{6.1b}
\end{equation*}
$$

For a free particle $\mathrm{V}(\mathrm{x})=$ constant, and then, the solutions to (6.1) are trivially

$$
\psi_{ \pm}(x)=\exp ( \pm i k x)=\exp ( \pm i p x / \hbar) .
$$

For a slowly varying potential function $\mathrm{V}(\mathrm{x})$ we now try a solution

$$
\begin{equation*}
\psi(x)=c_{+} \psi_{+}(x)+c_{-} \psi_{-}(x) . \tag{6.2}
\end{equation*}
$$

This is sc. Wentzel-Kramers-Brillouin (WKB) approximation.
Next, let us make a substitution

$$
\pm \mathrm{px}=\mathrm{S}_{ \pm}(\mathrm{x})=\mathrm{S}_{ \pm}^{(0)}+\hbar \mathrm{S}_{ \pm}^{(1)}(\mathrm{x})
$$

QTES, 2021
This leads to a solution (see the text book)

$$
\begin{equation*}
\psi(x)=\left(\frac{1}{p(x)}\right)^{1 / 2}\left\{c_{+} \exp \left[\frac{\mathrm{i}}{\hbar} \int_{0}^{x} p(x) \mathrm{d} x\right]+c_{-} \exp \left[-\frac{\mathrm{i}}{\hbar} \int_{0}^{x} p(x) \mathrm{d} x\right]\right\} \tag{6.5b}
\end{equation*}
$$

The probability density of $\psi(x)$ in Eq. (6.5b) is proportional to $1 / \mathrm{p}$. This is the same as the "classical probability density" 1/v.

For the classically allowed region, $\mathrm{E}>\mathrm{V}$, where p is real, the solution (6.5b) takes the form

$$
\begin{equation*}
\psi(x)=\frac{C}{p^{1 / 2}(x)} \sin \left\{\frac{1}{\hbar} \int_{0}^{x} p(x) \mathrm{d} x+\delta\right\} \tag{6.6}
\end{equation*}
$$

where $\mathrm{C} \sin \delta=\left(\mathrm{c}_{+}+\mathrm{c}_{-}\right)$. In this region WKB approximation is good, see the Fig. 6.1.

At $\mathrm{E}=\mathrm{V}$ WKB diverges at $\mathrm{p}=0$, and in classically forbidden region $p$ is imaginary. In the latter case the solution to (6.5) is a linear combination of exponentially decaying and exponentially growing functions.


Fig. 6.1.

## Time-independent perturbation theory

Consider a system, whose hamiltonian H is almost the same as the one of the known "reference system" $\mathrm{H}^{(0)}$. Let us denote the deviation or perturbation as $\mathrm{H}^{(1)}$. Thus,

$$
\begin{equation*}
\mathrm{H}=\mathrm{H}^{(0)}+\mathrm{H}^{(1)} . \tag{6.10}
\end{equation*}
$$

Here we assume that both of the hamiltonians are time-independent.

### 6.1. Perturbation of two-level system

Assume we know solutions of the reference system, i.e., solutions to the equation

$$
\begin{equation*}
\mathrm{H}^{(0)} \psi^{(0)}{ }_{\mathrm{m}}=\mathrm{E}^{(0)}{ }_{\mathrm{m}} \psi^{(0)}{ }_{\mathrm{m}} . \tag{6.11a}
\end{equation*}
$$

Let us further assume a reference system of two states, only: $\psi^{(0)}{ }_{m} ; m=1,2$; in bracket notation, normalized levels 11$\rangle$ and 12$\rangle$. Now, we search for solutions of the perturbed system

$$
\begin{equation*}
\mathrm{H} \psi=\mathrm{E} \psi \tag{6.11b}
\end{equation*}
$$

by using a trial wave function

$$
\begin{equation*}
\psi=\mathrm{a}_{1} \psi^{(0)}{ }_{1}+\mathrm{a}_{2} \psi^{(0)}{ }_{2}=\mathrm{a}_{1}|1\rangle+\mathrm{a}_{2}|2\rangle . \tag{6.12}
\end{equation*}
$$

Thus, we need to find $\mathrm{a}_{1}$ and $\mathrm{a}_{2}$.

QTES, 2021
Substitution of (6.12) to (6.11b) gives

$$
H\left(a_{1}|1\rangle+a_{2}|2\rangle\right)=E\left(a_{1}|1\rangle+a_{2}|2\rangle\right)
$$

and multiplication by both $\langle 11$ and $\langle 2|$ gives

$$
\begin{array}{lll}
\mathrm{a}_{1}\langle 1| \mathrm{H}|1\rangle+\mathrm{a}_{2}\langle 1| \mathrm{H}|2\rangle=E \mathrm{a}_{1} \\
\mathrm{a}_{1}\langle 2| \mathrm{H}|1\rangle+\mathrm{a}_{2}\langle 2| \mathrm{H}|2\rangle=\mathrm{E} \mathrm{a}_{2} & \text { or } & \mathrm{a}_{1} \mathrm{H}_{11}+\mathrm{a}_{2} \mathrm{H}_{12}=\mathrm{E} \mathrm{a} \mathrm{a}_{1} \\
\mathrm{a}_{1} \mathrm{H}_{21}+\mathrm{a}_{2} \mathrm{H}_{22}=\mathrm{E} \mathrm{a}
\end{array}
$$

as $\langle\mathrm{ilj}\rangle=\delta_{\mathrm{ij}}$. This is a pair of coupled equations for the unknown factors $\mathrm{a}_{1}$ and $\mathrm{a}_{2}$

$$
\begin{align*}
\left(H_{11}-E\right) a_{1}+H_{12} a_{2} & =0 \\
H_{21} & a_{1}+\left(H_{22}-E\right) a_{2} \tag{6.13}
\end{align*}=0, ~ \$
$$

which has non-trivial solutions only if

This implies, that

$$
\left|\begin{array}{cc}
\mathrm{H}_{11}-\mathrm{E} & \mathrm{H}_{12}  \tag{6.14a}\\
\mathrm{H}_{21} & \mathrm{H}_{22}-\mathrm{E}
\end{array}\right|=0 .
$$

$$
\begin{equation*}
\left(\mathrm{H}_{11}-\mathrm{E}\right)\left(\mathrm{H}_{22}-\mathrm{E}\right)-\mathrm{H}_{12} \mathrm{H}_{21}=0, \tag{6.14b}
\end{equation*}
$$

with solutions

$$
\begin{equation*}
\mathrm{E}_{ \pm}=1 / 2\left(\mathrm{H}_{11}+\mathrm{H}_{22}\right) \pm 1 / 2\left\{\left(\mathrm{H}_{22}-\mathrm{H}_{11}\right)^{2}+4 \mathrm{H}_{12} \mathrm{H}_{21}\right\}^{1 / 2} . \tag{6.14c}
\end{equation*}
$$

Now, consider a special case, where $\mathrm{H}_{\mathrm{mm}}=\mathrm{H}^{(0)}{ }_{\mathrm{mm}}+\mathrm{H}^{(1)}{ }_{\mathrm{mm}}=\mathrm{H}^{(0)}{ }_{\mathrm{mm}}=\mathrm{E}^{(0)}{ }_{\mathrm{m}}$, i.e., $\mathrm{H}^{(1)}{ }_{\mathrm{mm}}=0$.
As for the nondiagonals we have $\mathrm{H}_{12}=\mathrm{H}^{(0)}{ }_{12}+\mathrm{H}^{(1)}{ }_{12}=\mathrm{H}^{(1)}{ }_{12}$ and similarly $\mathrm{H}_{21}=\mathrm{H}^{(1)}{ }_{21}$, it follows

$$
\begin{equation*}
\mathrm{E}_{ \pm}=1 / 2\left(\mathrm{E}^{(0)}{ }_{1}+\mathrm{E}^{(0)}{ }_{2}\right) \pm 1 / 2\left\{\left(\mathrm{E}^{(0)}{ }_{2}-\mathrm{E}^{(0)}{ }_{1}\right)^{2}+4 \varepsilon^{2}\right\}^{1 / 2}, \tag{6.15}
\end{equation*}
$$

where $\varepsilon^{2}=\mathrm{H}^{(1)}{ }_{12} \mathrm{H}^{(1)}{ }_{21}$. Furthermore, $\varepsilon^{2}=\left|\mathrm{H}^{(1)}{ }_{12}\right|^{2}$, because H is hermitian.

Fig 6.4 shows how the "perturbation" $\varepsilon$ and $\Delta \mathrm{E}=\mathrm{E}^{(0)}{ }_{2}-\mathrm{E}^{(0)}{ }_{1}$ cause increasing separation of the levels $\mathrm{E}^{(0)}{ }_{1}$ and $\mathrm{E}^{(0)}{ }_{2}$.
If $\varepsilon / \Delta \mathrm{E} \ll 1$, by using $(1+\mathrm{x})^{1 / 2} \approx 1+1 / 2 \mathrm{x}$, where $\mathrm{x} \ll 1$, we obtain
and

$$
\begin{align*}
& \mathrm{E}_{+} \approx \mathrm{E}^{(0)}{ }_{2}+\varepsilon^{2} / \Delta \mathrm{E}^{(0)}  \tag{6.16}\\
& \mathrm{E}_{-} \approx \mathrm{E}^{(0)}{ }_{1}-\varepsilon^{2} / \Delta \mathrm{E}^{(0)} . \tag{6.17a}
\end{align*}
$$

Let us find wave functions by using trial functions
$\psi_{-}=\cos \beta \psi^{(0)}{ }_{1}+\sin \beta \psi^{(0)}{ }_{2}$ and $\psi_{+}=-\sin \beta \psi^{(0)}{ }_{1}+\cos \beta \psi^{(0)}{ }_{2}$.
These are "already orthonormalized", if $|1\rangle=\psi_{1}$ and $|2\rangle=$ $\psi_{2}$ are, because $\langle+\mid+\rangle=\langle-\mid-\rangle=\sin ^{2} \beta+\cos ^{2} \beta=1$ and

Figs. 6.3 and 6.4. $\langle+I-\rangle=\sin \beta \cos \beta-\sin \beta \cos \beta=0$. Let us find $\beta$ by substituting $I-\rangle=\psi$ - into the
Schrödinger equation and using the orthogonality condition

$$
0=\langle+| \mathrm{H}|-\rangle=-\sin \beta \cos \beta \mathrm{H}_{11}+\cos ^{2} \beta \mathrm{H}_{12}-\sin ^{2} \beta \mathrm{H}_{21}+\sin \beta \cos \beta \mathrm{H}_{22},
$$

which implies $\left(\mathrm{E}^{(0)}{ }_{1}-\mathrm{E}^{(0)}{ }_{2}\right) \sin \beta \cos \beta=\cos ^{2} \beta \mathrm{H}^{(1)}{ }_{12}-\sin ^{2} \beta \mathrm{H}^{(1)}{ }_{21}$. If $\mathrm{H}^{(1)}{ }_{12}=\mathrm{H}^{(1)}{ }_{21}$, it follows

$$
\begin{equation*}
\tan 2 \beta=2\left|\mathrm{H}^{(1)}{ }_{12}\right| /\left(\mathrm{E}^{(0)}{ }_{2}-\mathrm{E}^{(0)}{ }_{1}\right) . \tag{6.17b}
\end{equation*}
$$

In case of degeneracy, $\left(\mathrm{E}^{(0)}{ }_{1}-\mathrm{E}^{(0)}{ }_{2}\right)=0$, this leads to $\tan 2 \beta=\infty$ or $\sin \beta=\cos \beta=1 / \sqrt{ } 2$ and

$$
\psi_{-}=1 / \sqrt{ } 2\left(\psi_{1}+\psi_{2}\right) \quad \text { and } \quad \psi_{+}=1 / \sqrt{ } 2\left(-\psi_{1}+\psi_{2}\right) .
$$

In case of small perturbation, $\mathrm{H}^{(1)}{ }_{12} / \Delta \mathrm{E} \ll 1$, we have $\tan 2 \beta=2 \mathrm{H}^{(1)}{ }_{12} /\left(\mathrm{E}_{2}-\mathrm{E}_{1}\right) \approx 2 \beta \ll 1$, and $\sin \beta \approx \beta$ and $\cos \beta \approx 1$. Thus,

$$
\begin{equation*}
\underline{\psi_{-} \approx \psi^{(0)}{ }_{1}+\left(\mathrm{H}^{(1)} 12 / \Delta \mathrm{E}^{(0)}\right) \psi^{(0)}{ }_{2} \quad \text { and } \quad \psi_{+} \approx \psi^{(0)}{ }_{2}-\left(\mathrm{H}^{(1)}{ }_{12} / \Delta \mathrm{E}^{(0)}\right) \psi^{(0)}{ }_{1} .} \tag{6.18}
\end{equation*}
$$

QTES, 2021

### 6.2. Many-level systems

Consider a general case, where the (unperturbed) reference system is

$$
\begin{equation*}
\left.\mathrm{H}^{(0)}|\mathrm{n}\rangle=\mathrm{E}_{\mathrm{n}}{ }^{(0)} \ln \right\rangle ; \mathrm{n}=0,1, \ldots \tag{6.19}
\end{equation*}
$$

and number of states is not restricted, but we assume the states are not degenerate. Now, we write

$$
\mathrm{H}=\mathrm{H}^{(0)}+\mathrm{H}^{(1)}+\mathrm{H}^{(2)}+\ldots .
$$

Let us define $\lambda$ (strength of perturbation) to keep track of the order of perturbation, as

$$
\begin{equation*}
\mathrm{H}=\mathrm{H}^{(0)}+\lambda \mathrm{H}^{(1)}+\lambda^{2} \mathrm{H}^{(2)}+\ldots . \tag{6.20a}
\end{equation*}
$$

At the end we set $\lambda=1$. Similarly, we write the wave function
and eigenenergy

$$
\begin{gather*}
\psi=\psi_{0}{ }^{(0)}+\lambda \psi_{0}^{(1)}+\lambda^{2} \psi_{0}^{(2)}+\ldots  \tag{6.20b}\\
\mathrm{E}_{0}=\mathrm{E}_{0}{ }^{(0)}+\lambda \mathrm{E}_{0}^{(1)}+\lambda^{2} \mathrm{E}_{0}^{(2)}+\ldots .
\end{gather*}
$$

Substitution to

$$
\mathrm{H} \psi=\mathrm{E} \psi .
$$

gives

$$
\begin{aligned}
& \lambda^{0}\left\{\mathrm{H}^{(0)} \psi_{0}{ }_{0}^{(0)} \mathrm{E}_{0} \psi_{0}{ }^{(0)}\right\}+\lambda\left\{\mathrm{H}^{(0)} \psi_{0}{ }^{(1)}+\mathrm{H}^{(1)} \psi_{0}{ }^{(0)}-\mathrm{E}_{0}{ }^{(0)} \psi_{0}{ }^{(1)}-\mathrm{E}_{0}{ }^{(1)} \psi_{0}{ }^{(0)}\right\}+ \\
+ & \lambda^{2}\left\{\mathrm{H}^{(0)} \psi_{0}(2)+\mathrm{H}^{(1)} \psi_{0}{ }^{(1)}+\mathrm{H}^{(2)} \psi_{0}^{(0)}-\mathrm{E}_{0} \psi_{0}^{(2)}-\mathrm{E}^{(1)} 0 \psi_{0}^{(1)}-\mathrm{E}^{(2)} 0 \psi_{0}{ }^{(0)}\right\}+\ldots=0 .
\end{aligned}
$$

Considering $\lambda$ as arbitrary, each of the orders must vanish separately, and thus,

$$
\begin{gather*}
\mathrm{H}^{(0)} \psi_{0}{ }^{(0)}=\mathrm{E}_{0}{ }^{(0)} \psi_{0}{ }^{(0)}  \tag{6.21a}\\
\left(\mathrm{H}^{(0)}-\mathrm{E}_{0}{ }^{(0)}\right) \psi_{0}{ }^{(1)}=\left(\mathrm{E}_{0}{ }^{(1)}-\mathrm{H}^{(1)}\right) \psi_{0}{ }^{(0)}  \tag{6.21b}\\
\left(\mathrm{H}^{(0)}-\mathrm{E}_{0}^{(0)}\right) \psi_{0}{ }^{(2)}=\left(\mathrm{E}_{0}{ }^{(2)}-\mathrm{H}^{(2)}\right) \psi_{0}{ }^{(0)}+\left(\mathrm{E}_{0}{ }^{(1)}-\mathrm{H}^{(1)}\right) \psi_{0}^{(1)} . \tag{6.21c}
\end{gather*}
$$

## The first order correction to energy

Eq. (6.21a) is simply for the ground state $(\mathrm{n}=0$ ) of Eq. (6.19). Note that similar expansions are valid for other states ( $\mathrm{n}>0$ ), as well.
Let us search for the first order correction to the ground state $\psi_{0}{ }^{(0)}$ with a trial expansion

$$
\begin{equation*}
\left.\psi_{0}{ }^{(1)}=\sum_{\mathrm{n}} \mathrm{a}_{\mathrm{n}} \psi^{(0)}{ }_{\mathrm{n}}=\sum_{\mathrm{n}} \mathrm{a}_{\mathrm{n}} \ln \right\rangle, \tag{6.22}
\end{equation*}
$$

substituted to (6.21b). This gives

$$
\begin{equation*}
\left.\Sigma_{\mathrm{n}} \mathrm{a}_{\mathrm{n}}\left(\mathrm{E}_{\mathrm{n}}{ }^{(0)}-\mathrm{E}_{0}{ }^{(0)}\right) \ln \right\rangle=\left(\mathrm{E}_{0}^{(1)}-\mathrm{H}^{(1)}\right)|0\rangle \tag{6.23}
\end{equation*}
$$

Multiplying by $\left\langle 0\right.$ from the left we obtain $\left.0=\mathrm{E}_{0}{ }^{(1)}-\langle 0| \mathrm{H}^{(1)} \mid 0\right\rangle$, which gives the first order correction to the energy

$$
\begin{equation*}
\mathrm{E}_{0}{ }^{(1)}=\langle 0| \mathrm{H}^{(1)}|0\rangle . \tag{6.24}
\end{equation*}
$$

## The first order correction to wavefunction

Multiplying (6.23) from the left by $\langle\mathrm{kl}$, for $\mathrm{k} \neq 0$, we get

$$
\mathrm{a}_{\mathrm{k}}\left(\mathrm{E}^{(0)}{ }_{\mathrm{k}}-\mathrm{E}^{(0)}{ }_{0}\right)=-\mathrm{H}^{(1)}{ }_{\mathrm{k} 0}
$$

and

$$
\begin{equation*}
\mathrm{a}_{\mathrm{k}}=\mathrm{H}^{(1)}{ }_{\mathrm{k} 0} /\left(\mathrm{E}^{(0)}{ }_{0}-\mathrm{E}^{(0)}{ }_{\mathrm{k}}\right) . \tag{6.26}
\end{equation*}
$$

Thus, the first order wavefunction is

$$
\begin{equation*}
\psi_{0}=\psi_{0}{ }^{(0)}+\psi_{0}^{(1)}=\psi_{0}^{(0)}+\sum_{\mathrm{k} \neq 0}\left\{\mathrm{H}^{(1)}{ }_{\mathrm{k} 0} /\left(\mathrm{E}^{(0)}{ }_{0}-\mathrm{E}^{(0)} \mathrm{k}\right)\right\} \psi^{(0)}{ }_{\mathrm{k}} . \tag{6.27}
\end{equation*}
$$

QTES, 2021

## The second order correction to energy

Let us use the same technique for finding the second order correction to the energy, So,

$$
\begin{equation*}
\psi^{(2)}{ }_{0}=\sum_{\mathrm{n}} \mathrm{~b}_{\mathrm{n}} \psi^{(0)_{n}}=\sum_{\mathrm{n}} \mathrm{~b}_{\mathrm{n}}|\mathrm{n}\rangle, \tag{6.28}
\end{equation*}
$$

which we substitute into (6.21c). Similarly as above, we get for the second order correction

$$
\begin{equation*}
\mathrm{E}_{0}^{(2)}=\mathrm{H}_{00}^{(2)}+\sum_{\mathrm{n}}^{\prime}\left(\frac{\mathrm{H}_{0 \mathrm{n}}^{(1)} \mathrm{H}_{\mathrm{n} 0}^{(1)}}{\mathrm{E}_{0}-\mathrm{E}_{\mathrm{n}}}\right) . \tag{6.30}
\end{equation*}
$$

### 6.3.Comments on perturbation expressions

It can be shown, that the wavefunction of order n is sufficient to give the energy of order $2 \mathrm{n}+1$.
The main factor effecting on the convergence of the PT expansion is the strength of the perturbation. Another significant factor is the possible disappearance of matrix elements

$$
\begin{equation*}
\mathrm{H}^{(1)} 0 \mathrm{n}=\langle 0| \mathrm{H}^{(1)}|\mathrm{n}\rangle \tag{6.31}
\end{equation*}
$$

because of symmetry, i.e., as "vanishing integrals".

## The closure approximation

In case the denominator in (6.30) can be approximated as $\mathrm{E}_{0}-\mathrm{E}_{\mathrm{n}} \approx-\Delta \mathrm{E}$,

$$
\begin{equation*}
\mathrm{E}_{0}{ }^{(2)} \approx \mathrm{H}_{00}{ }^{(2)}-\left\{\sum_{\mathrm{n}} \mathrm{H}_{0 \mathrm{n}}{ }^{(1)} \mathrm{H}_{\mathrm{n} 0}{ }^{(1)}-\mathrm{H}_{00}{ }^{(1)} \mathrm{H}_{00}{ }^{(1)}\right\} / \Delta \mathrm{E} \tag{6.32}
\end{equation*}
$$

and denoting the numerator as

$$
\begin{equation*}
\varepsilon^{2}=\langle 0| \mathrm{H}^{(1) 2}|0\rangle-\langle 0| \mathrm{H}^{(1)}|0\rangle^{2}, \tag{6.33a}
\end{equation*}
$$

we can write

$$
\begin{equation*}
\mathrm{E}_{0}{ }^{(2)} \approx \mathrm{H}_{00}{ }^{(2)}-\varepsilon^{2} / \Delta \mathrm{E} . \tag{6.33b}
\end{equation*}
$$



Fig. 6.7.

### 6.4. Perturbation theory for degenerate states

The above perturbation theory does not apply for degenerate states, because of divergent denominators in the expressions of wavefunctions and energies. A modification is needed.
Consider r-fold degenerate eigenenergy $\mathrm{E}_{0}$ from the reference system equation

$$
\begin{equation*}
\mathrm{H}^{(0)}|0, \ell\rangle=\mathrm{E}^{(0)}{ }_{0}|0, \ell\rangle ; \ell=1,2, \ldots, \mathrm{r} \tag{6.35}
\end{equation*}
$$

where $|0, \ell\rangle$ are the linearly independent states. And as before, we have $H=H^{(0)}+H^{(1)}$. It is useful to choose such linear combinations of $I 0, \ell\rangle$ which fit to the symmetry of the perturbation. Therefore, we write

$$
\begin{equation*}
\left.\varphi^{(0)}{ }_{0 \mathrm{i}}=\sum_{\ell} \mathrm{d}_{\ell \mathrm{i}} \mathrm{I} 0, \ell\right\rangle=\sum_{\ell} \mathrm{d}_{\ell \mathrm{i}} \psi^{(0)}{ }_{\ell \mathrm{i}} \tag{6.36}
\end{equation*}
$$

which should diagonalize $\mathrm{H}^{(1)}$, i.e., $\left\langle\varphi^{(0)}{ }_{0 i}\right| \mathrm{H}^{(1)}\left|\varphi^{(0)}{ }_{0 \mathrm{j}}\right\rangle=0$, if $\mathrm{i} \neq \mathrm{j}$. Now, let us find the coefficients $\mathrm{d}_{\ell \mathrm{i}}$ doing that.
As before, we start with
and

$$
\begin{aligned}
& \varphi_{i}=\psi^{(0)}{ }_{0 i}+\lambda \psi^{(1)}{ }_{0 i}+\ldots \\
& E_{i}=E^{(0)}{ }_{0 i}+\lambda E^{(1)}{ }_{0 i}+\ldots,
\end{aligned}
$$

substitute these into $H \psi_{i}=E_{i} \psi_{i}$ and obtain

$$
\begin{align*}
\mathrm{H}^{(0)} \varphi^{(0)} 0 \mathrm{i} & =\mathrm{E}^{(0)} 0 \varphi^{(0)} 0 \mathrm{i}  \tag{6.37a}\\
\left(\mathrm{H}^{(0)}-\mathrm{E}^{(0)} 0\right) \psi^{(1)}{ }_{0 \mathrm{i}} & =\left(\mathrm{E}_{0}(1){ }_{0 \mathrm{i}}-\mathrm{H}^{(1)}\right) \varphi^{(0)}{ }_{0 \mathrm{i}}, \tag{6.37b}
\end{align*}
$$

as Eqs. (6.21) above.

Now, let us write the first order correction as

$$
\begin{equation*}
\psi^{(1)} 0 \mathrm{i}=\sum_{\ell} \mathrm{c}_{\ell} \psi^{(0)}{ }_{0 \ell}+\sum_{\mathrm{n} \neq 0} \mathrm{c}_{\mathrm{n}} \psi^{(0)_{\mathrm{n}}} \tag{6.38}
\end{equation*}
$$

where the former sum is over the degenerate states $\psi^{(0)}{ }_{0 \ell}$ and the latter is over all others. With notation $\left.\psi^{(0)}{ }_{\ell}=10 \ell\right\rangle$ and substitution with Eq. (6.36) $\varphi^{(0)}{ }_{0 i}=\sum_{\mathrm{r}} \mathrm{d}_{\ell \mathrm{i}}|0, \ell\rangle$ into (6.37b) we get

$$
\begin{equation*}
\left.\sum_{\ell} \mathrm{c}_{\ell}\left(\mathrm{E}^{(0)}{ }_{0}-\mathrm{E}^{(0)}{ }_{0}\right)|0 \ell\rangle+\sum_{\mathrm{n} \neq 0} \mathrm{c}_{\mathrm{n}}\left(\mathrm{E}^{(0)}{ }_{\mathrm{n}}-\mathrm{E}^{(0)}{ }_{0}\right) \ln \right\rangle=\sum_{\ell} \mathrm{d}_{\ell \mathrm{i}}\left(\mathrm{E}^{(1)}{ }_{0 \mathrm{i}}-\mathrm{H}^{(1)}\right)|0 \ell\rangle . \tag{6.39a}
\end{equation*}
$$

Multiplying this from left by $\langle 0 \mathrm{kl}$, which is one (possibly linear combination) of functions $\langle 0 \ell$ l such that $\langle 0 \mathrm{k} \mid \mathrm{n}\rangle=0$ and $\langle 0 \mathrm{k} \mid 0 \ell\rangle=\mathrm{S}_{\mathrm{k} \ell} \neq 0$, we get sc. secular equations

$$
\begin{equation*}
0=\sum_{\ell} \mathrm{d}_{\ell \mathrm{i}}\left\{\mathrm{E}^{(1)}{ }_{0 \mathrm{i}} \mathrm{~S}_{\mathrm{k} \ell}-\langle 0 \mathrm{k}| \mathrm{H}^{(1)}|0 \ell\rangle\right\} \tag{6.39b}
\end{equation*}
$$

or

$$
\begin{equation*}
\sum_{\ell} \mathrm{d}_{\ell \mathrm{i}}\left(\mathrm{E}^{(1)}{ }_{0 \mathrm{i}} \mathrm{~S}_{\mathrm{k} \ell}-\mathrm{H}^{(1)}{ }_{\mathrm{k} \ell}\right)=0, \tag{6.41}
\end{equation*}
$$

where
and

$$
\begin{equation*}
\mathrm{S}_{\mathrm{k} \ell}=\langle 0 \mathrm{k} \mid 0 \ell\rangle \tag{6.40a}
\end{equation*}
$$

This group of equations, $\mathrm{i}=1,2, \ldots, \mathrm{r}$; (or matrix equation) for coefficients $\mathrm{d}_{\ell \mathrm{i}}(\ell=1,2, \ldots, \mathrm{r})$, has nontrivial solutions, if the secular determinant vanishes, i.e.,

$$
\begin{equation*}
\operatorname{det}\left[\mathrm{H}^{(1)} \mathrm{k}_{\ell}-\mathrm{E}^{(1)}{ }_{0 \mathrm{i}} \mathrm{~S}_{\mathrm{k} \ell}\right]_{\mathrm{i} \ell}=0 . \tag{6.42}
\end{equation*}
$$

The first order energy corrections $\mathrm{E}^{(1)}{ }_{0 \mathrm{i}}(\mathrm{i}=1,2, \ldots, \mathrm{r})$ are found from matrix diagonalization. The, the coefficients $\mathrm{d}_{l i}$ corresponding to each $\mathrm{E}^{(1)}{ }_{0 \mathrm{i}}$ are found from Eq. (6.41).
Note, if we choose $\mathrm{S}_{\mathrm{k} \ell}=\delta_{\mathrm{k} \ell}$ and $\mathrm{r}=2$, the Eq. (6.42) returns the earlier two-level case, Eq. (6.14c).

## Variation theory

Another approach to find or estimate the solution to the Schrödinger equation is to use a trial wavefunction $\psi_{\text {trial }}$. It can be written by using a proper functional form (educated guess) with parameters, which will be fitted or "optimized" to make $\psi_{\text {trial }}$ the best approximate to the exact wavefunction.

### 6.5. Variation theorem

Define Rayleigh ratio

$$
\begin{equation*}
\mathcal{E}=\left\langle\psi_{\text {trial }}\right| \mathrm{H}\left|\psi_{\text {trial }}\right\rangle /\left\langle\psi_{\text {trial }} \mid \psi_{\text {trial }}\right\rangle . \tag{6.43a}
\end{equation*}
$$

Then, we can use the variation theorem

$$
\begin{equation*}
\mathcal{E} \geq \mathrm{E}_{0} \quad \text { for any } \psi_{\text {trial }} \tag{6.43b}
\end{equation*}
$$

where $\mathrm{E}_{0}$ is the ground state energy of the hamiltonian H . The equality $\mathcal{E}=\mathrm{E}_{0}$ holds only, if the the trial function is identical with the exact wavefunction, $\psi_{\text {trial }}=\psi_{0}$.
Let us prove the variation theorem by writing $\psi_{\text {trial }}=\sum_{n} c_{n} \psi_{n}=\sum_{n} c_{n}|n\rangle$, where $\left\{\psi_{n}\right\}$ is the complete set of solutions to $\mathrm{H} \psi_{\mathrm{n}}=\mathrm{E}_{\mathrm{n}} \psi_{\mathrm{n}}$. Because

$$
\begin{aligned}
\left\langle\psi_{\text {trial }}\right|\left(\mathrm{H}-\mathrm{E}_{0}\right)\left|\psi_{\text {trial }}\right\rangle & =\sum_{\mathrm{n}, \mathrm{~m}} \mathrm{c}_{\mathrm{n}} * \mathrm{c}_{\mathrm{m}}\langle\mathrm{n}|\left(\mathrm{H}-\mathrm{E}_{0}\right)|\mathrm{m}\rangle \\
& =\sum_{\mathrm{n}, \mathrm{~m}} \mathrm{c}_{\mathrm{n}}^{*} \mathrm{c}_{\mathrm{m}}\left(\mathrm{E}_{\mathrm{m}}-\mathrm{E}_{0}\right)\langle\mathrm{n} \mid \mathrm{m}\rangle=\sum_{\mathrm{n}}\left|\mathrm{c}_{\mathrm{n}}\right|^{2}\left(\mathrm{E}_{\mathrm{n}}-\mathrm{E}_{0}\right) \geq 0,
\end{aligned}
$$

we obtain $\left\langle\psi_{\text {trial }}\right| \mathrm{H}\left|\psi_{\text {trial }}\right\rangle \geq \mathrm{E}_{0}\left\langle\psi_{\text {trial }} \mid \psi_{\text {trial }}\right\rangle$, which proves the theorem.
Optimization of the parameters $\left\{p_{n}\right\}$ is done the "usual way" from the conditions

$$
\left(\partial \mathcal{E} / \partial p_{1}\right)=0,\left(\partial \mathcal{E} / \partial p_{2}\right)=0, \ldots ;
$$

i.e., from the extremum of the gradient of $\mathcal{E}$ in the parameter space $\left\{\mathrm{p}_{\mathrm{n}}\right\}$.

Example 6.8. Use the trial $\psi_{\text {trial }}(\mathrm{r})=\mathrm{e}^{-\mathrm{kr}}$ for the ground state of hydrogen like one-electron atoms with nuclear charge Ze ). Find the parameter k and the corresponding eigenenergy.

$$
\int_{0}^{\infty} \mathrm{r}^{\mathrm{n}} \mathrm{e}^{-\alpha \mathrm{r}} \mathrm{dr}=\frac{\mathrm{n}!}{\alpha^{\mathrm{n}+1}}
$$

### 6.6. Rayleigh-Ritz method

The Rayleigh-Ritz method is based on a trial wavefunction

$$
\begin{equation*}
\psi_{\text {trial }}=\sum_{\mathrm{i}} \mathrm{c}_{\mathrm{i}} \psi_{\mathrm{i}}=\sum_{\mathrm{i}} \mathrm{c}_{\mathrm{i}}|\mathrm{i}\rangle \tag{6.44}
\end{equation*}
$$

where the parameters $c_{i}$ are to be optimized. Functions $\left\{\psi_{i}\right\}_{i=0}{ }^{\mathrm{N}-1}$ are called as basis set. Now, the Rayleighin ratio is
$\mathcal{E}=\left\langle\psi_{\text {trial }} \mid \mathrm{Hl} \psi_{\text {trial }}\right\rangle /\left\langle\psi_{\text {trial }} \mid \psi_{\text {trial }}\right\rangle=\left\{\sum_{\mathrm{ij}} \mathrm{c}_{\mathrm{i}} * \mathrm{c}_{\mathrm{j}}\langle\mathrm{ilHlj}\rangle\right\} /\left\{\sum_{\mathrm{ij}} \mathrm{c}_{\mathrm{i}}{ }^{*} \mathrm{c}_{\mathrm{j}}\langle\mathrm{ilj}\rangle\right\}=\left\{\Sigma_{\mathrm{ij}} \mathrm{c}_{\mathrm{i}} * \mathrm{c}_{\mathrm{j}} \mathrm{H}_{\mathrm{ij}}\right\} /\left\{\sum_{\mathrm{ij}} \mathrm{c}_{\mathrm{i}} * \mathrm{c}_{\mathrm{j}} \mathrm{S}_{\mathrm{ij}}\right\}$. Allowing real coefficients $c_{i}$, only, minimization condition for $\mathcal{E}$ is

This implies, that for all k

$$
\Sigma_{\mathrm{i}} \mathrm{c}_{\mathrm{i}}\left(\mathrm{H}_{\mathrm{ik}}-\mathcal{E} \mathrm{S}_{\mathrm{ik}}\right)=0
$$

and thus, non-trivial solutions exist, if

$$
\operatorname{det}\left(\mathrm{H}_{\mathrm{ik}}-\mathcal{E} \mathrm{S}_{\mathrm{ik}}\right)=0
$$

$$
\begin{align*}
& \quad \begin{array}{l}
\frac{d \mathcal{E}}{d c_{k}}= \\
\left(\sum_{\mathrm{j}} \mathrm{c}_{\mathrm{j}} H_{\mathrm{kj}}+\sum_{\mathrm{i}} \mathrm{c}_{\mathrm{i}} H_{\mathrm{ik}}\right) \sum_{\mathrm{ij}} \mathrm{c}_{\mathrm{i}} \mathrm{c}_{\mathrm{j}} \mathrm{~S}_{\mathrm{ij}}-\left(\sum_{\mathrm{j}} \mathrm{c}_{\mathrm{j}} \mathrm{~S}_{\mathrm{kj}}+\sum_{\mathrm{i}} \mathrm{c}_{\mathrm{i}} \mathrm{~S}_{\mathrm{ik}}\right) \sum_{\mathrm{ij}} \mathrm{c}_{\mathrm{i}} \mathrm{c}_{\mathrm{j}} H_{\mathrm{ij}} \\
\left(\sum_{\mathrm{ij}} \mathrm{c}_{\mathrm{i}} \mathrm{c}_{\mathrm{j}} \mathrm{~S}_{\mathrm{ij}}\right)^{2}
\end{array}= \\
& \quad=\frac{\sum_{\mathrm{j}} \mathrm{c}_{\mathrm{j}}\left(\mathrm{H}_{\mathrm{kj}}-E \mathrm{~S}_{\mathrm{kj}}\right)+\sum_{\mathrm{i}} \mathrm{c}_{\mathrm{i}}\left(\mathrm{H}_{\mathrm{ik}}-E \mathrm{~S}_{\mathrm{ik}}\right)}{\sum_{\mathrm{ij}} \mathrm{c}_{\mathrm{i}} \mathrm{c}_{\mathrm{j}} \mathrm{~S}_{\mathrm{ij}}}=0 .
\end{align*}
$$

This leads to N equations and
N roots, i.e., N eigenenergies $\mathcal{F}_{\mathrm{i}}$, where N is the number of basis functions in (6.44). The lowest one, $\mathfrak{E}_{0}$, is the ground state energy. For all energies the coefficients are obtained from (6.46) and the wavefunction from (6.44).

QTES, 2021100

## Hellmann-Feynman theorem

For a system of quantum particles (atoms, molecules, nanostrutures) the structure and external interactions are described in the hamiltonian. In the following we inspect how the total energy of the system is changed, when the hamiltonian is changed.
Consider hamiltonian with a parameter P (e.g. a bond length or an external field, etc.). Then, both the wavefunction $\psi$ and the energy $\mathrm{E}(\mathrm{P})=\langle\psi| \mathrm{H}|\psi\rangle$ depend on P . Assume the wavefunction is normalized $\langle\psi \mid \psi\rangle=1$. Now, the change of energy as a function of $P$ is

$$
\begin{aligned}
& \qquad \begin{aligned}
\mathrm{dE} / \mathrm{dP} & =\mathrm{d} / \mathrm{dP}\langle\psi| \mathrm{H}|\psi\rangle=\langle\mathrm{d} \psi / \mathrm{dP}| \mathrm{H}|\psi\rangle+\langle\psi| \mathrm{dH} / \mathrm{dP}|\psi\rangle+\langle\psi| \mathrm{H}|\mathrm{~d} \psi / \mathrm{dP}\rangle \\
& =\mathrm{E} \mathrm{~d} / \mathrm{dP}\langle\psi \mid \psi\rangle+\langle\psi| \mathrm{dH} / \mathrm{dP}|\psi\rangle=\langle\mathrm{dH} / \mathrm{dP}\rangle .
\end{aligned} \\
& \text { Thus, we have proved the Hellmann-Feynman theorem: } \quad \left\lvert\, \begin{array}{l}
\text { H. Hellmann, Einführung in die } \\
\text { Quantenchemie (1937). } \\
\text { R.P. Feynman, PR56, 340 (1939). }
\end{array}\right.
\end{aligned}
$$

$$
\begin{equation*}
\frac{\mathrm{dE}}{\mathrm{dP}}=\left\langle\frac{\mathrm{dH}}{\mathrm{dP}}\right\rangle \tag{6.48}
\end{equation*}
$$

This is a practical method for calculation of forces in molecular dynamics simulations. For practical accuracy, the wavefunction should be known accurately, enough.
Example Find the molecular energy as a function of the external electric field $\mathcal{E}$, when $\mathrm{H}=\mathrm{H}^{(0)}-\mu_{\mathrm{z}} \mathcal{E}$

## Time-dependent perturbation theory

### 6.7. Time-dependence of two-level system

Time-dependent perturbation theory is needed to consider the system response to turning the perturbation "on" or "off", and also, if the perturbation operator explicitly depends on time,

$$
\begin{equation*}
\mathrm{H}(\mathrm{t})=\mathrm{H}^{(0)}+\mathrm{H}^{(1)}(\mathrm{t}) . \tag{6.49}
\end{equation*}
$$

A usual time-dependent perturbation is the electromagnetic field

$$
\begin{equation*}
\mathrm{H}^{(1)}(\mathrm{t})=\mathrm{H}^{(1)} 2 \cos (\omega \mathrm{t}) . \tag{6.50}
\end{equation*}
$$

Let us find the time-dependent solution $\Psi(\mathrm{t})$ of the Schrödinger equation

$$
\begin{equation*}
\mathrm{H} \Psi=\mathrm{i} \hbar(\partial \Psi / \partial \mathrm{t}) \tag{6.51}
\end{equation*}
$$

Let us first restrict us to a two-level system (as we did in sec. 6.1), where the eigenenergies $\mathrm{E}_{1}$ and $\mathrm{E}_{2}$ correspond to eigenstates $\psi_{1}$ and $\psi_{2}$ of the stationary state reference system

$$
\begin{equation*}
\mathrm{H}^{(0)} \psi_{\mathrm{n}}=\mathrm{E}_{\mathrm{n}} \psi_{\mathrm{n}} ; \mathrm{n}=1,2 ; \quad \text { with } \quad \Psi_{\mathrm{n}}(\mathrm{t})=\psi_{\mathrm{n}} \mathrm{e}^{-\mathrm{i} \mathrm{E}_{\mathrm{n}} \mathrm{t} / \hbar} \tag{6.52}
\end{equation*}
$$

from Eq. (1.38). Let us choose a trial wavefunction

$$
\begin{equation*}
\Psi(\mathrm{t})=\mathrm{c}_{1}(\mathrm{t}) \Psi_{1}(\mathrm{t})+\mathrm{c}_{2}(\mathrm{t}) \Psi_{2}(\mathrm{t}) \tag{6.53}
\end{equation*}
$$

where also the coefficients $\mathrm{c}_{\mathrm{i}}(\mathrm{t})$ depend on time. With substitution to (6.51) we obtain $\mathrm{c}_{1} \mathrm{H}^{(0)} \Psi_{1}+\mathrm{c}_{1} \mathrm{H}^{(1)} \Psi_{1}+\mathrm{c}_{2} \mathrm{H}^{(0)} \Psi_{2}+\mathrm{c}_{2} \mathrm{H}^{(1)} \Psi_{2}=\mathrm{i} \hbar\left(\partial \mathrm{c}_{1} / \partial \mathrm{t} \Psi_{1}+\mathrm{c}_{1} \partial \Psi_{1} / \partial \mathrm{t}+\partial \mathrm{c}_{2} / \partial \mathrm{t} \Psi_{2}+\mathrm{c}_{2} \partial \Psi_{2} / \partial \mathrm{t}\right)$ and by using

$$
\mathrm{H}^{(0)} \Psi_{\mathrm{n}}=\mathrm{i} \hbar\left(\partial \Psi_{\mathrm{n}} / \partial \mathrm{t}\right)
$$

we get

$$
\begin{equation*}
\mathrm{c}_{1} \mathrm{H}^{(1)}(\mathrm{t}) \Psi_{1}+\mathrm{c}_{2} \mathrm{H}^{(1)}(\mathrm{t}) \Psi_{2}=\mathrm{i} \hbar \dot{\mathrm{c}}_{1} \Psi_{1}+\mathrm{i} \hbar \dot{\mathrm{c}}_{2} \Psi_{2}, \tag{6.54}
\end{equation*}
$$

where we denote $\dot{c}=\mathrm{dc} / \mathrm{dt}$.

QTES, 2021102
With substitution of (6.52), we get

$$
\begin{equation*}
\mathrm{c}_{1} \mathrm{H}^{(1)}(\mathrm{t}) \psi_{1} \mathrm{e}^{-\mathrm{i} \mathrm{E}_{1} t / \hbar}+\mathrm{c}_{2} \mathrm{H}^{(1)}(\mathrm{t}) \psi_{2} \mathrm{e}^{-\mathrm{i} \mathrm{E}_{2} t / \hbar}=\mathrm{i} \hbar \dot{\mathrm{c}}_{1} \psi_{1} \mathrm{e}^{-\mathrm{i} \mathrm{E}_{1} t / \hbar}+\mathrm{i} \hbar \dot{\mathrm{c}}_{2} \psi_{2} \mathrm{e}^{-\mathrm{i} \mathrm{E}_{2} t / \hbar} \tag{6.55a}
\end{equation*}
$$

and further, multiplying by $\psi_{1} *$ and integrating $\left(\int \psi_{1} * \psi_{2} d \tau=0\right)$, we get

$$
\begin{equation*}
\mathrm{c}_{1} \mathrm{H}^{(1)}{ }_{11}(\mathrm{t}) \mathrm{e}^{-\mathrm{iE} \mathrm{E}_{1} \mathrm{t} / \hbar}+\mathrm{c}_{2} \mathrm{H}^{(1)}{ }_{12}(\mathrm{t}) \mathrm{e}^{-\mathrm{i} \mathrm{E}_{2} \mathrm{t} / \hbar}=\mathrm{i} \hbar \dot{\mathrm{c}}_{1} \mathrm{e}^{-\mathrm{i} \mathrm{E}_{1} \mathrm{t} / \hbar} \tag{6.55b}
\end{equation*}
$$

where

$$
\mathrm{H}^{(1)}{ }_{\mathrm{ij}}(\mathrm{t})=\int \psi_{\mathrm{i}}^{*} \mathrm{H}^{(1)}(\mathrm{t}) \psi_{\mathrm{j}} \mathrm{~d} \tau=\left\langle\psi_{\mathrm{i}}\right| \mathrm{H}^{(1)}(\mathrm{t})\left|\psi_{\mathrm{j}}\right\rangle .
$$

Now, let us write for the difference $\mathrm{E}_{2}-\mathrm{E}_{1}=\hbar \omega_{21}$ and assume that $\mathrm{H}^{(1)}{ }_{11}(\mathrm{t})=\mathrm{H}^{(1)}{ }_{22}(\mathrm{t})=0$, again. This is valid for the usual perturbations, like electromagnetic field. Then, we find

$$
\begin{align*}
& \dot{\mathrm{c}}_{1}=(1 / \mathrm{i} \hbar) \mathrm{c}_{2} \mathrm{H}^{(1)}{ }_{12}(\mathrm{t}) \mathrm{e}^{-\mathrm{i} \omega_{21} \mathrm{t}}  \tag{6.57a}\\
& \dot{\mathrm{c}}_{2}=(1 / \mathrm{i} \hbar) \mathrm{c}_{1} \mathrm{H}^{(1)}{ }_{21}(\mathrm{t}) \mathrm{e}^{\mathrm{i} \omega_{21} \mathrm{t}} . \tag{6.57b}
\end{align*}
$$

and similarly

Now, let us discuss two cases:
(i) The perturbation is "off" or absent. Then $\dot{\mathrm{c}}_{1}=\dot{\mathrm{c}}_{2}=0$,
$c_{1}=c_{1}(0)=$ constant, $c_{2}=c_{2}(0)=$ constant and

$$
\begin{equation*}
\Psi(\mathrm{t})=\mathrm{c}_{1}(0) \psi_{1} \mathrm{e}^{-\mathrm{iE}} \mathrm{E}_{1} \mathrm{t} / \hbar \quad+\mathrm{c}_{2}(0) \psi_{2} \mathrm{e}^{-\mathrm{iE}} \mathrm{E}_{2} \mathrm{t} / \hbar \tag{6.58}
\end{equation*}
$$

Now, $\left|c_{1}(0) e^{-i E_{1} t / \hbar}\right|^{2}=\left|c_{1}(0)\right|^{2}$ is a constant, whose interpretation is the probability that the system occupies the state $\psi_{1}$. The other constant, $\left|c_{2}(0)\right|^{2}$, is the probability for occupation of the state $\psi_{2}$.
(ii) Assume the perturbation is "on" with a constant strength, i.e., $\mathrm{H}^{(1)}{ }_{12}(\mathrm{t})=\hbar \mathrm{V}$ and $\mathrm{H}^{(1)} 21(\mathrm{t})=\hbar \mathrm{V}^{*}$ are constants. Then, solution to the pair of equations (6.57)

$$
\begin{aligned}
& \dot{c}_{1}=(1 / i \hbar) c_{2} \mathrm{H}^{(1)}{ }_{12} \mathrm{e}^{-\mathrm{i} \omega_{21} \mathrm{t}} \\
& \dot{\mathrm{c}}_{2}=(1 / i \hbar) \mathrm{c}_{1} \mathrm{H}^{(1)}{ }_{21} \mathrm{e}^{\mathrm{i} \omega_{21} \mathrm{t}} .
\end{aligned}
$$

will be found from

$$
\begin{align*}
\ddot{\mathrm{c}}_{2} & =(1 / \mathrm{i} \hbar) \dot{\mathrm{c}}_{1} \mathrm{H}^{(1)}{ }_{21} \mathrm{e}^{\mathrm{i} \omega_{21} \mathrm{t}}+\mathrm{i} \omega_{21}(1 / \mathrm{i} \hbar) \mathrm{c}_{1} \mathrm{H}^{(1)}{ }_{21} \mathrm{e}^{\mathrm{i} \omega_{21} \mathrm{t}} \\
& =(1 / \mathrm{i} \hbar)^{2} \mathrm{c}_{2} \mathrm{H}^{(1)}{ }_{12} \mathrm{H}^{(1)}{ }_{21}+\mathrm{i} \omega_{21} \dot{\mathrm{c}}_{2} \\
& =-\left.\mathrm{IV}\right|^{2} \mathrm{c}_{2}+\mathrm{i} \omega_{21} \dot{\mathrm{c}}_{2}, \tag{6.60}
\end{align*}
$$

with the notation $\mathrm{H}^{(1)}{ }_{12} \mathrm{H}^{(1)}{ }_{21}=\hbar^{2} \mathrm{~V}^{2}$. Solution to this is

$$
\begin{equation*}
c_{2}(t)=\left(A e^{i \Omega t}+B e^{-i \Omega t}\right) e^{i \omega_{21} t / 2} ; \Omega=1 / 2\left(\omega_{21}^{2}+4|V|^{2}\right)^{1 / 2}, \tag{6.61}
\end{equation*}
$$

where $A$ and $B$ are given by the initial conditions. Similar expression is found for $c_{1}$, and if the initial conditions are $\mathrm{c}_{1}(0)=1$ and $\mathrm{c}_{2}(0)=0$, then

$$
\begin{align*}
& \mathrm{c}_{1}(\mathrm{t})=\left\{\cos (\Omega \mathrm{t})+\mathrm{i}\left(\omega_{21} / 2 \Omega\right) \sin (\Omega \mathrm{t})\right\} \mathrm{e}^{-\mathrm{i} \omega_{21} \mathrm{t} / 2} \\
& \mathrm{c}_{2}(\mathrm{t})=-\mathrm{i}(|\mathrm{~V}| / \Omega) \sin (\Omega \mathrm{t}) \mathrm{e}^{\mathrm{i} \omega_{21} t / 2} . \tag{6.62}
\end{align*}
$$

and
These are exact solutions in case of two levels.

## Rabi oscillations

Consider the probabilities of finding the system in state 1 or 2 . Let us denote these as $\mathrm{P}_{1}$ and $\mathrm{P}_{2}$. $\mathrm{P}_{1}+\mathrm{P}_{2}=1$, as there are no other states. Thus, $\mathrm{P}_{2}(\mathrm{t})=1-\mathrm{P}_{1}(\mathrm{t})$ with the time dependence called Rabi oscillations

$$
\begin{equation*}
P_{2}=\left.\dot{I}_{c_{2}}\right|^{2}=\frac{4 \mathrm{~V}^{2}}{\omega_{21}^{2}+4 \mathrm{~V}^{2}} \sin ^{2} \frac{\left(\omega_{21}^{2}+4 \mathrm{~V}^{2}\right)^{1 / 2}}{2} t \tag{6.63}
\end{equation*}
$$

Now, let us consider two common cases:
(i) A degenerate system, where $\mathrm{E}_{1}=\mathrm{E}_{2}$, and thus, $\omega_{21}=0$. Then

$$
\begin{equation*}
P_{2}(t)=\sin ^{2}(I V \mid t), \tag{6.64}
\end{equation*}
$$

shown in Fig. 6.12. We see, that the larger the perturbation, the faster the oscillation; but on the other hand, no matter how weak the perturbation is, it is always sufficient for the "transition".
(ii) The other extreme $\left(E_{2}-E_{1}\right) / \hbar \gg V$, whence

$$
\begin{equation*}
\mathrm{P}_{2}(\mathrm{t}) \approx\left(2|\mathrm{~V}| / \omega_{21}\right)^{2} \sin ^{2}\left(1 / 2 \omega_{21} \mathrm{t}\right) \tag{6.65}
\end{equation*}
$$

This is illustrated in Fig. 6.13. Now, we see that oscillation frequency is given by the level spacing and the amplitude by the perturbation strength with respect to level spacing. Occupation probability of level 2 is always lower than 1.

### 6.8. Many-level systems

The order of the differential equations grows together with the number of levels, and therefore, general solutions cannot be found similarly to the two-level case.
For book keeping of the "virtual transitions" between the levels in many-level case, one can use the sc. Feynman diagrams.

### 6.9. Fermi's golden rule

Consider a transititon in a many-level system from the initial state $\mathrm{E}_{\mathrm{i}}$ to the final state $\mathrm{E}_{\mathrm{f}}$. Denote $\mathrm{E}_{\mathrm{f}}-\mathrm{E}_{\mathrm{i}}=\hbar \omega_{\mathrm{fi}}$. Assume interaction with the electromagnetic radiation

$$
\begin{equation*}
\mathrm{H}^{(1)}(\mathrm{t})=\mathrm{H}^{(1)} 2 \cos \omega \mathrm{t}=\mathrm{H}^{(1)}\left(\mathrm{e}^{\mathrm{i} \omega \mathrm{t}}+\mathrm{e}^{-\mathrm{i} \omega \mathrm{t}}\right) . \tag{6.75}
\end{equation*}
$$

Then, it can be shown that occupation probability of the final state $f$, if $i$ is initially unoccupied, is

$$
\begin{equation*}
P_{\mathrm{fi}}(\mathrm{t})=\frac{4 \mathrm{~V}_{\mathrm{fi}}^{2}}{\left(\omega_{\mathrm{fi}}-\omega\right)^{2}} \sin ^{2} \frac{\omega_{\mathrm{fi}}-\omega}{2} t \tag{6.77}
\end{equation*}
$$

where $\mathrm{V}_{\mathrm{fi}}{ }^{2}=\mathrm{H}^{(1)}{ }_{\text {if }} \mathrm{H}^{(1)}{ }_{\mathrm{fi}} / \hbar_{2}=\mathrm{V}_{\mathrm{if}}{ }^{2}=\mathrm{V}^{2}$. Compare with the twolevel system, Eq. (6.65). We see, that occupation probability of state f strongly increases while approaching the resonance $\omega_{\mathrm{fi}}-\omega=0$.
In case of several closely spaced final states we denote the density of states (DOS) by $\rho\left(\mathrm{E}_{\mathrm{f}}\right)$, at about $\mathrm{E}_{\mathrm{f}}$. Then, it can be shown that

$$
\begin{equation*}
\mathrm{P}_{\mathrm{fi}}(\mathrm{t})=2 \pi \hbar \mathrm{~V}_{\mathrm{fi}}^{2} \rho\left(\mathrm{E}_{\mathrm{f}}\right) \mathrm{t} \tag{6.82}
\end{equation*}
$$

and by defining the transition rate (spectral line intensity) as

$$
\begin{equation*}
\mathrm{W}_{\mathrm{if}}=\mathrm{dP}_{\mathrm{fi}} / \mathrm{dt} \tag{6.83}
\end{equation*}
$$

we arrive at the Fermi's golden rule

$$
\begin{equation*}
\mathrm{W}_{\mathrm{if}}=2 \pi \hbar \mathrm{~V}_{\mathrm{fi}}^{2} \rho\left(\mathrm{E}_{\mathrm{f}}\right)=2 \pi / \hbar \mid \mathrm{H}^{(1)_{\mathrm{fi}}} \mathrm{I}^{2} \rho\left(\mathrm{E}_{\mathrm{f}}\right) . \tag{6.84}
\end{equation*}
$$

QTES, 2021106

### 6.10. Einstein transition probabilities ( $A$ and $B$ coefficients)

Transition probability is proportional also to the intensity of the electromagnetic field $\rho(v)$ (number of photons / units of time and space), at the transition frequency. Thus, the rate (or probability) of the stimulated absorption can be written as

$$
\begin{equation*}
W_{\text {if }}=B_{\text {if }} \rho, \tag{6.86}
\end{equation*}
$$

where $\mathrm{B}_{\text {if }}$ is sc. Einstein $B$ coefficient for stimulated absorption. Correspondingly, the transition rate (or probability) of stimulated emission can be written as

$$
\begin{equation*}
\mathrm{W}_{\mathrm{fi}}=\mathrm{B}_{\mathrm{fi}} \rho, \tag{6.87}
\end{equation*}
$$

where $B_{f i}$ is sc. Einstein $B$ coefficient for stimulated emission. As $\mathrm{IH}^{(1)} \mathrm{fi}^{2}=\left.\mathrm{IH}^{(1)}{ }_{\mathrm{if}}\right|^{2}, \mathrm{~B}_{\mathrm{if}}=\mathrm{B}_{\mathrm{fi}}$.
In thermal equilibrium the occupation of energy levels obey the Boltzmann distribution, i.e., $\mathrm{N}_{\mathrm{f}} / \mathrm{N}_{\mathrm{i}}=\exp (-\mathrm{h} v / \mathrm{kT})$. Now, because $\mathrm{W}_{\mathrm{if}}=\mathrm{W}_{\mathrm{fi}}$, for an equilibrium, there has to be one more transition mechanism. This is the spontaneous emission, for which

$$
\begin{equation*}
\mathrm{W}_{\mathrm{fi}}=\mathrm{A}_{\mathrm{fi}} . \tag{6.90}
\end{equation*}
$$

In equilibrium $\mathrm{N}_{\mathrm{i}} \mathrm{W}_{\mathrm{if}}=\mathrm{N}_{\mathrm{f}} \mathrm{W}_{\mathrm{fi}}$, whence $\mathrm{N}_{\mathrm{i}} \mathrm{B} \rho=\mathrm{N}_{\mathrm{f}}(\mathrm{B} \rho+\mathrm{A})$,
where $\mathrm{A}=\mathrm{A}_{\mathrm{fi}}$ and $\mathrm{B}=\mathrm{B}_{\mathrm{if}}=\mathrm{B}_{\mathrm{fi}}$. From comparison with the Boltzmann distribution we get

$$
\mathrm{N}_{\mathrm{f}} / \mathrm{N}_{\mathrm{i}}=\mathrm{B} \rho /(\mathrm{B} \rho+\mathrm{A})=\exp (-\mathrm{h} v / \mathrm{kT}),
$$

and further

$$
\begin{equation*}
\rho=(\mathrm{A} / \mathrm{B}) /\{\exp (\mathrm{h} v / \mathrm{kT})-1\} . \tag{6.92a}
\end{equation*}
$$

Now, comparison with the Planck distribution of black body radiation, Eq. (0.5)

$$
\begin{equation*}
\rho(v)=\left(8 \pi \mathrm{~h} v^{3} / \mathrm{c}^{3}\right) /\{\exp (\mathrm{h} v / \mathrm{kT})-1\} \tag{6.92b}
\end{equation*}
$$

gives

$$
\begin{equation*}
\mathrm{A} / \mathrm{B}=8 \pi \mathrm{~h}(\mathrm{v} / \mathrm{c})^{3} . \tag{6.93}
\end{equation*}
$$

Thus, the proportion of spontaneous emission to stimulated emission increases as the cube of transition frequency (or level spacing).

### 6.11. Lifetime and spectral linewidth (energy uncertainty)

Some unspecified time after excitation, the excited state will decay to some lower state, e.g., the ground state, through spontaneous emission. All such transitions contribute to the lifetime of the excited state.
The stationary state with an eigenenergy E , is described by the wavefunction (6.52) $\Psi(\mathrm{t})=\psi \mathrm{e}^{-\mathrm{iEt} / \hbar}$ and $|\Psi|^{2}=|\psi|^{2}$, but for a decaying state

$$
\begin{equation*}
|\Psi|^{2}=|\psi|^{2} e^{-t / \tau} \tag{6.94}
\end{equation*}
$$

and therefore,

$$
\begin{equation*}
\Psi=\psi \mathrm{e}^{-\mathrm{iEt} / h-t / 2 \tau} \tag{6.95}
\end{equation*}
$$

where $\tau$ is the lifetime of the state.

QTES, 2021
This kind of decaying wavefunction is a superposition of "several" wavefuctions with energies close to E . This superposition wavefunction is

$$
\begin{equation*}
\mathrm{e}^{-\mathrm{E} \mathrm{E} t h-1 / 2 \tau}=\int \mathrm{g}\left(\mathrm{E}^{\prime}\right) \mathrm{e}^{-\mathrm{i} \cdot} \mathrm{E}^{\prime} / \hbar \mathrm{dE}^{\prime}, \tag{6.96a}
\end{equation*}
$$

i.e., a Fourier transform of the spectral line (spectral density function)

$$
\begin{equation*}
\mathrm{g}\left(\mathrm{E}^{\prime}\right)=\frac{\hbar / \tau}{\left(\mathrm{E}-\mathrm{E}^{\prime}\right)^{2}+(\hbar / 2 \tau)^{2}} \tag{6.96b}
\end{equation*}
$$

This is sc. Lorentz line shape of spontaneous emission of photons carrying away the energy of transition, i.e., $\Delta \mathrm{E}=\mathrm{E}_{\mathrm{i}}-\mathrm{E}_{\mathrm{f}}=\hbar \omega_{\mathrm{if}}$.

Thus, the function $g\left(E^{\prime}\right)$ describes the energy spectrum of contributing frequencies or wavelengths. So, if the state is metastable, its energy is not discrete, but a distribution with Lorentzian shape and a half width $\delta \mathrm{E}$, an uncertainty of the excited state energy, broadening of the spectral line.
Fourier analysis gives for the half width $\delta \mathrm{E}=\hbar / \tau$.
Thus, it is customary to present this as an
"uncertainty relation of time and energy"

$$
\begin{equation*}
\tau \delta \mathrm{E} \approx \hbar . \tag{6.97}
\end{equation*}
$$

This can be used to determine the lifetimes of quantum states from experimental line widths.

