By ir. J.C.A. Wevers



Dear reader,
This document contains a 108 page LATEX file which contains a lot equations in physics. It is written at advanced undergraduate/postgraduate level. It is intended to be a short reference for anyone who works with physics and often needs to look up equations.

This, and a Dutch version of this file, can be obtained from the author, Johan Wevers (johanw@vulcan.xs4all.nl).

It can also be obtained on the WWW. See http://www.xs4all.nl/~johanw/index.html, where also a Postscript version is available.

If you find any errors or have any comments, please let me know. I am always open for suggestions and possible corrections to the physics formulary.

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Johan Wevers

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## Physical Constants

| Name | Symbol | Value | Unit |
| :---: | :---: | :---: | :---: |
| Number $\pi$ | $\pi$ | 3.1415926535897932 |  |
| Number e | $\begin{aligned} & \mathrm{e} \\ & \gamma=\lim _{n \rightarrow \infty}\left(\sum_{k=1}^{n} 1 / k-\ln (n)\right)=0.5772156649 \end{aligned}$ |  |  |
| Euler's constant |  |  |  |
| Elementary charge Gravitational constant Fine-structure constant Speed of light in vacuum Permittivity of the vacuum Permeability of the vacuum $\left(4 \pi \varepsilon_{0}\right)^{-1}$ | $\begin{aligned} & e \\ & G, \kappa \\ & \alpha=e^{2} / 2 h c \varepsilon_{0} \\ & c \\ & \varepsilon_{0} \\ & \mu_{0} \end{aligned}$ | $1.60217733 \cdot 10^{-19}$ | $\begin{aligned} & \mathrm{C} \\ & \mathrm{~m}^{3} \mathrm{~kg}^{-1} \mathrm{~s}^{-2} \\ & \mathrm{~m} / \mathrm{s}(\text { def }) \\ & \mathrm{F} / \mathrm{m} \\ & \mathrm{H} / \mathrm{m} \\ & \mathrm{Nm}^{2} \mathrm{C}^{-2} \\ & \hline \end{aligned}$ |
|  |  | $6.67259 \cdot 10^{-11}$ |  |
|  |  | $\approx 1 / 137$ |  |
|  |  | $2.99792458 \cdot 10^{8}$ |  |
|  |  | $8.854187 \cdot 10^{-12}$ |  |
|  |  | $4 \pi \cdot 10^{-7}$ |  |
|  |  | $8.9876 \cdot 10^{9}$ |  |
| Planck's constant <br> Dirac's constant <br> Bohr magneton <br> Bohr radius <br> Rydberg's constant <br> Electron Compton wavelength <br> Proton Compton wavelength <br> Reduced mass of the H -atom | $h$$\begin{aligned} & \hbar=h / 2 \pi \\ & \mu_{\mathrm{B}}=e \hbar / 2 m_{\mathrm{e}} \\ & a_{0} \\ & R y \\ & \lambda_{\mathrm{Ce}}=h / m_{\mathrm{e}} c \\ & \lambda_{\mathrm{Cp}}=h / m_{\mathrm{p}} c \\ & \mu_{\mathrm{H}} \\ & \hline \end{aligned}$ | $6.6260755 \cdot 10^{-34}$ | $\begin{aligned} & \hline \mathrm{Js} \\ & \mathrm{Js} \\ & \mathrm{Am}^{2} \\ & \AA \\ & \mathrm{eV} \\ & \mathrm{~m} \\ & \mathrm{~m} \\ & \mathrm{~kg} \\ & \hline \end{aligned}$ |
|  |  | $1.0545727 \cdot 10^{-34}$ |  |
|  |  | $9.2741 \cdot 10^{-24}$ |  |
|  |  | 0.52918 |  |
|  |  | 13.595 |  |
|  |  | $2.2463 \cdot 10^{-12}$ |  |
|  |  | $1.3214 \cdot 10^{-15}$ |  |
|  |  | $9.1045755 \cdot 10^{-31}$ |  |
| Stefan-Boltzmann's constant <br> Wien's constant | $\begin{aligned} & \sigma \\ & k_{\mathrm{W}} \end{aligned}$ | $5.67032 \cdot 10^{-8}$ | $\begin{aligned} & \mathrm{Wm}^{-2} \mathrm{~K}^{-4} \\ & \mathrm{mK} \end{aligned}$ |
|  |  | $2.8978 \cdot 10^{-3}$ |  |
| Molar gasconstant | $\begin{aligned} & \hline R \\ & N_{\mathrm{A}} \\ & k=R / N_{\mathrm{A}} \\ & \hline \end{aligned}$ | 8.31441 | $\begin{aligned} & \mathrm{J} \cdot \mathrm{~mol}^{-1} \cdot \mathrm{~K}^{-1} \\ & \mathrm{~mol}^{-1} \\ & \mathrm{~J} / \mathrm{K} \\ & \hline \end{aligned}$ |
| Avogadro's constant |  | $6.0221367 \cdot 10^{23}$ |  |
| Boltzmann's constant |  | $1.380658 \cdot 10^{-23}$ |  |
| Electron mass | $\begin{aligned} & m_{\mathrm{e}} \\ & m_{\mathrm{p}} \\ & m_{\mathrm{n}} \\ & m_{\mathrm{u}}=\frac{1}{12} m\left({ }_{6}^{12} \mathrm{C}\right) \\ & \mu_{\mathrm{N}} \\ & \hline \end{aligned}$ | $\begin{aligned} & \hline 9.1093897 \cdot 10^{-31} \\ & 1.6726231 \cdot 10^{-27} \\ & 1.674954 \cdot 10^{-27} \\ & 1.6605656 \cdot 10^{-27} \\ & 5.0508 \cdot 10^{-27} \end{aligned}$ | $\begin{aligned} & \mathrm{kg} \\ & \mathrm{~kg} \\ & \mathrm{~kg} \\ & \mathrm{~kg} \\ & \mathrm{~J} / \mathrm{T} \end{aligned}$ |
| Proton mass |  |  |  |
| Neutron mass |  |  |  |
| Elementary mass unit |  |  |  |
| Nuclear magneton |  |  |  |
| Diameter of the Sun | $D_{\odot}$ | $1392 \cdot 10^{6}$ | m |
| Mass of the Sun | $M_{\odot}$ | $1.989 \cdot 10^{30}$ | kg |
| Rotational period of the Sun | $T_{\text {¢ }}$ | 25.38 | days |
| Radius of Earth | $R_{\text {A }}$ | $6.378 \cdot 10^{6}$ | m |
| Mass of Earth | $M_{\text {A }}$ | $5.976 \cdot 10^{24}$ | kg |
| Rotational period of Earth | $T_{\text {A }}$ | 23.96 | hoursdays |
| Earth orbital period | Tropical year AU | 365.24219879 |  |
| Astronomical unit |  | $1.4959787066 \cdot 10^{11}$ | days m |
| Light year | 1j | $9.4605 \cdot 10^{15}$ | m |
| Parsec | pc | $3.0857 \cdot 10^{16}$ | $\begin{aligned} & \mathrm{m} \\ & \mathrm{~km} \cdot \mathrm{~s}^{-1} \cdot \mathrm{Mpc}^{-1} \\ & \hline \end{aligned}$ |
| Hubble constant | H | $\approx(75 \pm 25)$ |  |

## Chapter 1

## Mechanics

### 1.1 Point-kinetics in a fixed coordinate system

### 1.1.1 Definitions

The position $\vec{r}$, the velocity $\vec{v}$ and the acceleration $\vec{a}$ are defined by: $\vec{r}=(x, y, z), \vec{v}=(\dot{x}, \dot{y}, \dot{z}), \vec{a}=(\ddot{x}, \ddot{y}, \ddot{z})$. The following holds:

$$
s(t)=s_{0}+\int|\vec{v}(t)| d t ; \quad \vec{r}(t)=\vec{r}_{0}+\int \vec{v}(t) d t ; \quad \vec{v}(t)=\vec{v}_{0}+\int \vec{a}(t) d t
$$

When the acceleration is constant this gives: $v(t)=v_{0}+a t$ and $s(t)=s_{0}+v_{0} t+\frac{1}{2} a t^{2}$.
For the unit vectors in a direction $\perp$ to the orbit $\vec{e}_{\mathrm{t}}$ and parallel to it $\vec{e}_{\mathrm{n}}$ holds:

$$
\vec{e}_{\mathrm{t}}=\frac{\vec{v}}{|\vec{v}|}=\frac{d \vec{r}}{d s} \quad \dot{\vec{e}_{\mathrm{t}}}=\frac{v}{\rho} \vec{e}_{\mathrm{n}} ; \quad \vec{e}_{\mathrm{n}}=\frac{\dot{\overrightarrow{e_{\mathrm{t}}}}}{\left|\overrightarrow{\vec{e}_{\mathrm{t}}}\right|}
$$

For the curvature $k$ and the radius of curvature $\rho$ holds:

$$
\vec{k}=\frac{d \vec{e}_{\mathrm{t}}}{d s}=\frac{d^{2} \vec{r}}{d s^{2}}=\left|\frac{d \varphi}{d s}\right| ; \quad \rho=\frac{1}{|k|}
$$

### 1.1.2 Polar coordinates

Polar coordinates are defined by: $x=r \cos (\theta), y=r \sin (\theta)$. So, for the unit coordinate vectors holds: $\dot{\vec{e}}_{r}=\dot{\theta} \vec{e}_{\theta}, \dot{\vec{e}}_{\theta}=-\dot{\theta} \vec{e}_{r}$
The velocity and the acceleration are derived from: $\vec{r}=r \vec{e}_{r}, \vec{v}=\dot{r} \vec{e}_{r}+r \dot{\theta} \vec{e}_{\theta}, \vec{a}=\left(\ddot{r}-r \dot{\theta}^{2}\right) \vec{e}_{r}+(2 \dot{r} \dot{\theta}+r \ddot{\theta}) \vec{e}_{\theta}$.

### 1.2 Relative motion

For the motion of a point D w.r.t. a point Q holds: $\vec{r}_{\mathrm{D}}=\vec{r}_{\mathrm{Q}}+\frac{\vec{\omega} \times \vec{v}_{\mathrm{Q}}}{\omega^{2}}$ with $\overrightarrow{\mathrm{QD}}=\vec{r}_{\mathrm{D}}-\vec{r}_{\mathrm{Q}}$ and $\omega=\dot{\theta}$.
Further holds: $\alpha=\ddot{\theta} .{ }^{\prime}$ means that the quantity is defined in a moving system of coordinates. In a moving system holds:
$\vec{v}=\vec{v}_{\mathrm{Q}}+\vec{v}^{\prime}+\vec{\omega} \times \vec{r}^{\prime}$ and $\vec{a}=\vec{a}_{\mathrm{Q}}+\vec{a}^{\prime}+\vec{\alpha} \times \vec{r}^{\prime}+2 \vec{\omega} \times \vec{v}^{\prime}+\vec{\omega} \times\left(\vec{\omega} \times \vec{r}^{\prime}\right)$
with $\vec{\omega} \times\left(\vec{\omega} \times \vec{r}^{\prime}\right)=-\omega^{2} \vec{r}_{n}^{\prime}$

### 1.3 Point-dynamics in a fixed coordinate system

### 1.3.1 Force, (angular)momentum and energy

Newton's 2nd law connects the force on an object and the resulting acceleration of the object where the momentum is given by $\vec{p}=m \vec{v}$ :

$$
\vec{F}(\vec{r}, \vec{v}, t)=\frac{d \vec{p}}{d t}=\frac{d(m \vec{v})}{d t}=m \frac{d \vec{v}}{d t}+\vec{v} \frac{d m}{d t} \stackrel{m=\text { const }}{=} m \vec{a}
$$

Newton's 3rd law is given by: $\vec{F}_{\text {action }}=-\vec{F}_{\text {reaction }}$.
For the power $P$ holds: $P=\dot{W}=\vec{F} \cdot \vec{v}$. For the total energy $W$, the kinetic energy $T$ and the potential energy $U$ holds: $W=T+U ; \quad \dot{T}=-\dot{U}$ with $T=\frac{1}{2} m v^{2}$.
The kick $\vec{S}$ is given by: $\vec{S}=\Delta \vec{p}=\int \vec{F} d t$
The work $A$, delivered by a force, is $A=\int_{1}^{2} \vec{F} \cdot d \vec{s}=\int_{1}^{2} F \cos (\alpha) d s$
The torque $\vec{\tau}$ is related to the angular momentum $\vec{L}: \vec{\tau}=\dot{\vec{L}}=\vec{r} \times \vec{F}$; and
$\vec{L}=\vec{r} \times \vec{p}=m \vec{v} \times \vec{r},|\vec{L}|=m r^{2} \omega$. The following equation is valid:

$$
\tau=-\frac{\partial U}{\partial \theta}
$$

Hence, the conditions for a mechanical equilibrium are: $\sum \vec{F}_{i}=0$ and $\sum \vec{\tau}_{i}=0$.
The force of friction is usually proportional to the force perpendicular to the surface, except when the motion starts, when a threshold has to be overcome: $F_{\text {fric }}=f \cdot F_{\text {norm }} \cdot \vec{e}_{\mathrm{t}}$.

### 1.3.2 Conservative force fields

A conservative force can be written as the gradient of a potential: $\vec{F}_{\mathrm{cons}}=-\vec{\nabla} U$. From this follows that $\nabla \times \vec{F}=\overrightarrow{0}$. For such a force field also holds:

$$
\oint \vec{F} \cdot d \vec{s}=0 \Rightarrow U=U_{0}-\int_{r_{0}}^{r_{1}} \vec{F} \cdot d \vec{s}
$$

So the work delivered by a conservative force field depends not on the trajectory covered but only on the starting and ending points of the motion.

### 1.3.3 Gravitation

The Newtonian law of gravitation is (in GRT one also uses $\kappa$ instead of $G$ ):

$$
\vec{F}_{\mathrm{g}}=-G \frac{m_{1} m_{2}}{r^{2}} \vec{e}_{r}
$$

The gravitational potential is then given by $V=-G m / r$. From Gauss law it then follows: $\nabla^{2} V=4 \pi G \varrho$.

### 1.3.4 Orbital equations

If $V=V(r)$ one can derive from the equations of Lagrange for $\phi$ the conservation of angular momentum:

$$
\frac{\partial \mathcal{L}}{\partial \phi}=\frac{\partial V}{\partial \phi}=0 \Rightarrow \frac{d}{d t}\left(m r^{2} \phi\right)=0 \Rightarrow L_{z}=m r^{2} \phi=\mathrm{constant}
$$

For the radial position as a function of time can be found that:

$$
\left(\frac{d r}{d t}\right)^{2}=\frac{2(W-V)}{m}-\frac{L^{2}}{m^{2} r^{2}}
$$

The angular equation is then:

$$
\phi-\phi_{0}=\int_{0}^{r}\left[\frac{m r^{2}}{L} \sqrt{\frac{2(W-V)}{m}-\frac{L^{2}}{m^{2} r^{2}}}\right]^{-1} d r \stackrel{r^{-2} \text { field }}{=} \arccos \left(1+\frac{\frac{1}{r}-\frac{1}{r_{0}}}{\frac{1}{r_{0}}+k m / L_{z}^{2}}\right)
$$

If $F=F(r): L=$ constant, if $F$ is conservative: $W=$ constant, if $\vec{F} \perp \vec{v}$ then $\Delta T=0$ and $U=0$.

## Kepler's orbital equations

In a force field $F=k r^{-2}$, the orbits are conic sections with the origin of the force in one of the foci (Kepler's 1st law). The equation of the orbit is:

$$
r(\theta)=\frac{\ell}{1+\varepsilon \cos \left(\theta-\theta_{0}\right)}, \text { or: } x^{2}+y^{2}=(\ell-\varepsilon x)^{2}
$$

with

$$
\ell=\frac{L^{2}}{G \mu^{2} M_{\mathrm{tot}}} ; \quad \varepsilon^{2}=1+\frac{2 W L^{2}}{G^{2} \mu^{3} M_{\mathrm{tot}}^{2}}=1-\frac{\ell}{a} ; \quad a=\frac{\ell}{1-\varepsilon^{2}}=\frac{k}{2 W}
$$

$a$ is half the length of the long axis of the elliptical orbit in case the orbit is closed. Half the length of the short axis is $b=\sqrt{a \ell} . \varepsilon$ is the excentricity of the orbit. Orbits with an equal $\varepsilon$ are of equal shape. Now, 5 types of orbits are possible:

1. $k<0$ and $\varepsilon=0$ : a circle.
2. $k<0$ and $0<\varepsilon<1$ : an ellipse.
3. $k<0$ and $\varepsilon=1$ : a parabole.
4. $k<0$ and $\varepsilon>1$ : a hyperbole, curved towards the centre of force.
5. $k>0$ and $\varepsilon>1$ : a hyperbole, curved away from the centre of force.

Other combinations are not possible: the total energy in a repulsive force field is always positive so $\varepsilon>1$.
If the surface between the orbit covered between $t_{1}$ and $t_{2}$ and the focus C around which the planet moves is $A\left(t_{1}, t_{2}\right)$, Kepler's 2 nd law is

$$
A\left(t_{1}, t_{2}\right)=\frac{L_{\mathrm{C}}}{2 m}\left(t_{2}-t_{1}\right)
$$

Kepler's 3rd law is, with $T$ the period and $M_{\text {tot }}$ the total mass of the system:

$$
\frac{T^{2}}{a^{3}}=\frac{4 \pi^{2}}{G M_{\mathrm{tot}}}
$$

### 1.3.5 The virial theorem

The virial theorem for one particle is:

$$
\langle m \vec{v} \cdot \vec{r}\rangle=0 \Rightarrow\langle T\rangle=-\frac{1}{2}\langle\vec{F} \cdot \vec{r}\rangle=\frac{1}{2}\left\langle r \frac{d U}{d r}\right\rangle=\frac{1}{2} n\langle U\rangle \text { if } U=-\frac{k}{r^{n}}
$$

The virial theorem for a collection of particles is:

$$
\langle T\rangle=-\frac{1}{2}\left\langle\sum_{\text {particles }} \vec{F}_{i} \cdot \vec{r}_{i}+\sum_{\text {pairs }} \vec{F}_{i j} \cdot \vec{r}_{i j}\right\rangle
$$

These propositions can also be written as: $2 E_{\text {kin }}+E_{\text {pot }}=0$.

### 1.4 Point dynamics in a moving coordinate system

### 1.4.1 Apparent forces

The total force in a moving coordinate system can be found by subtracting the apparent forces from the forces working in the reference frame: $\vec{F}^{\prime}=\vec{F}-\vec{F}_{\text {app }}$. The different apparent forces are given by:

1. Transformation of the origin: $F_{\text {or }}=-m \vec{a}_{a}$
2. Rotation: $\vec{F}_{\alpha}=-m \vec{\alpha} \times \vec{r}^{\prime}$
3. Coriolis force: $F_{\text {cor }}=-2 m \vec{\omega} \times \vec{v}$
4. Centrifugal force: $\vec{F}_{\mathrm{cf}}=m \omega^{2} \vec{r}_{n}{ }^{\prime}=-\vec{F}_{\mathrm{cp}} ; \vec{F}_{\mathrm{cp}}=-\frac{m v^{2}}{r} \vec{e}_{r}$

### 1.4.2 Tensor notation

Transformation of the Newtonian equations of motion to $x^{\alpha}=x^{\alpha}(x)$ gives:

$$
\frac{d x^{\alpha}}{d t}=\frac{\partial x^{\alpha}}{\partial \bar{x}^{\beta}} \frac{d \bar{x}^{\beta}}{d t}
$$

The chain rule gives:

$$
\frac{d}{d t} \frac{d x^{\alpha}}{d t}=\frac{d^{2} x^{\alpha}}{d t^{2}}=\frac{d}{d t}\left(\frac{\partial x^{\alpha}}{\partial \bar{x}^{\beta}} \frac{d \bar{x}^{\beta}}{d t}\right)=\frac{\partial x^{\alpha}}{\partial \bar{x}^{\beta}} \frac{d^{2} \bar{x}^{\beta}}{d t^{2}}+\frac{d \bar{x}^{\beta}}{d t} \frac{d}{d t}\left(\frac{\partial x^{\alpha}}{\partial \bar{x}^{\beta}}\right)
$$

so:

$$
\frac{d}{d t} \frac{\partial x^{\alpha}}{\partial \bar{x}^{\beta}}=\frac{\partial}{\partial \bar{x}^{\gamma}} \frac{\partial x^{\alpha}}{\partial \bar{x}^{\beta}} \frac{d \bar{x}^{\gamma}}{d t}=\frac{\partial^{2} x^{\alpha}}{\partial \bar{x}^{\beta} \partial \bar{x}^{\gamma}} \frac{d \bar{x}^{\gamma}}{d t}
$$

This leads to:

$$
\frac{d^{2} x^{\alpha}}{d t^{2}}=\frac{\partial x^{\alpha}}{\partial \bar{x}^{\beta}} \frac{d^{2} \bar{x}^{\beta}}{d t^{2}}+\frac{\partial^{2} x^{\alpha}}{\partial \bar{x}^{\beta} \partial \bar{x}^{\gamma}} \frac{d \bar{x}^{\gamma}}{d t}\left(\frac{d \bar{x}^{\beta}}{d t}\right)
$$

Hence the Newtonian equation of motion

$$
m \frac{d^{2} x^{\alpha}}{d t^{2}}=F^{\alpha}
$$

will be transformed into:

$$
m\left\{\frac{d^{2} x^{\alpha}}{d t^{2}}+\Gamma_{\beta \gamma}^{\alpha} \frac{d x^{\beta}}{d t} \frac{d x^{\gamma}}{d t}\right\}=F^{\alpha}
$$

The apparent forces are taken from he origin to the effect side in the way $\Gamma_{\beta \gamma}^{\alpha} \frac{d x^{\beta}}{d t} \frac{d x^{\gamma}}{d t}$.

### 1.5 Dynamics of masspoint collections

### 1.5.1 The centre of mass

The velocity w.r.t. the centre of mass $\vec{R}$ is given by $\vec{v}-\vec{R}$. The coordinates of the centre of mass are given by:

$$
\vec{r}_{\mathrm{m}}=\frac{\sum m_{i} \vec{r}_{i}}{\sum m_{i}}
$$

In a 2-particle system, the coordinates of the centre of mass are given by:

$$
\vec{R}=\frac{m_{1} \vec{r}_{1}+m_{2} \vec{r}_{2}}{m_{1}+m_{2}}
$$

With $\vec{r}=\vec{r}_{1}-\vec{r}_{2}$, the kinetic energy becomes: $T=\frac{1}{2} M_{\mathrm{tot}} \dot{R}^{2}+\frac{1}{2} \mu \dot{r}^{2}$, with the reduced mass $\mu$ given by: $\frac{1}{\mu}=\frac{1}{m_{1}}+\frac{1}{m_{2}}$
The motion within and outside the centre of mass can be separated:

$$
\begin{aligned}
& \dot{\vec{L}}_{\text {outside }}=\vec{\tau}_{\text {outside }} ; \quad \dot{\vec{L}}_{\text {inside }}=\vec{\tau}_{\text {inside }} \\
& \vec{p}=m \vec{v}_{\mathrm{m}} ; \quad \vec{F}_{\text {ext }}=m \vec{a}_{\mathrm{m}} ; \quad \vec{F}_{12}=\mu \vec{u}
\end{aligned}
$$

### 1.5.2 Collisions

With collisions, where B are the coordinates of the collision and C an arbitrary other position, holds: $\vec{p}=m \vec{v}_{\mathrm{m}}$ is constant, and $T=\frac{1}{2} m \vec{v}_{\mathrm{m}}^{2}$ is constant. The changes in the relative velocities can be derived from: $\vec{S}=\Delta \vec{p}=$ $\mu\left(\vec{v}_{\text {aft }}-\vec{v}_{\text {before }}\right)$. Further holds $\Delta \vec{L}_{\mathrm{C}}=\overrightarrow{\mathrm{CB}} \times \vec{S}, \vec{p} \| \vec{S}=$ constant and $\vec{L}$ w.r.t. B is constant.

### 1.6 Dynamics of rigid bodies

### 1.6.1 Moment of Inertia

The angular momentum in a moving coordinate system is given by:

$$
\vec{L}^{\prime}=I \vec{\omega}+\vec{L}_{n}^{\prime}
$$

where $I$ is the moment of inertia with respect to a central axis, which is given by:

$$
I=\sum_{i} m_{i} \vec{r}_{i}^{2} ; \quad T^{\prime}=W_{\text {rot }}=\frac{1}{2} \omega I_{i j} \vec{e}_{i} \vec{e}_{j}=\frac{1}{2} I \omega^{2}
$$

or, in the continuous case:

$$
I=\frac{m}{V} \int{r^{\prime}}_{n}^{2} d V=\int{r^{\prime}}_{n}^{2} d m
$$

Further holds:

$$
L_{i}=I^{i j} \omega_{j} ; \quad I_{i i}=I_{i} ; \quad I_{i j}=I_{j i}=-\sum_{k} m_{k} x_{i}^{\prime} x_{j}^{\prime}
$$

Steiner's theorem is: $I_{\text {w.r.t.D }}=I_{\text {w.r.t.C }}+m(D M)^{2}$ if axis $\mathrm{C} \|$ axis D.

| Object | $I$ | Object | $I$ |
| :--- | :--- | :--- | :--- |
| Cavern cylinder | $I=m R^{2}$ | Massive cylinder | $I=\frac{1}{2} m R^{2}$ |
| Disc, axis in plane disc through m | $I=\frac{1}{4} m R^{2}$ | Halter | $I=\frac{1}{2} \mu R^{2}$ |
| Cavern sphere | $I=\frac{2}{3} m R^{2}$ | Massive sphere | $I=\frac{2}{5} m R^{2}$ |
| Bar, axis $\perp$ through c.o.m. | $I=\frac{1}{12} m l^{2}$ | Bar, axis $\perp$ through end | $I=\frac{1}{3} m l^{2}$ |
| Rectangle, axis $\perp$ plane thr. c.o.m. | $I=\frac{1}{12} m\left(a^{2}+b^{2}\right)$ | Rectangle, axis $\\| b$ thr. m | $I=m a^{2}$ |

### 1.6.2 Principal axes

Each rigid body has (at least) 3 principal axes which stand $\perp$ to each other. For a principal axis holds:

$$
\frac{\partial I}{\partial \omega_{x}}=\frac{\partial I}{\partial \omega_{y}}=\frac{\partial I}{\partial \omega_{z}}=0 \text { so } L_{n}^{\prime}=0
$$

The following holds: $\dot{\omega}_{k}=-a_{i j k} \omega_{i} \omega_{j}$ with $a_{i j k}=\frac{I_{i}-I_{j}}{I_{k}}$ if $I_{1} \leq I_{2} \leq I_{3}$.

### 1.6.3 Time dependence

For torque of force $\vec{\tau}$ holds:

$$
\vec{\tau}^{\prime}=I \ddot{\theta} ; \quad \frac{d^{\prime \prime} \vec{L}^{\prime}}{d t}=\vec{\tau}^{\prime}-\vec{\omega} \times \vec{L}^{\prime}
$$

The torque $\vec{T}$ is defined by: $\vec{T}=\vec{F} \times \vec{d}$.

### 1.7 Variational Calculus, Hamilton and Lagrange mechanics

### 1.7.1 Variational Calculus

Starting with:

$$
\delta \int_{a}^{b} \mathcal{L}(q, \dot{q}, t) d t=0 \text { with } \delta(a)=\delta(b)=0 \text { and } \delta\left(\frac{d u}{d x}\right)=\frac{d}{d x}(\delta u)
$$

the equations of Lagrange can be derived:

$$
\frac{d}{d t} \frac{\partial \mathcal{L}}{\partial \dot{q}_{i}}=\frac{\partial \mathcal{L}}{\partial q_{i}}
$$

When there are additional conditions applying to the variational problem $\delta J(u)=0$ of the type $K(u)=$ constant, the new problem becomes: $\delta J(u)-\lambda \delta K(u)=0$.

### 1.7.2 Hamilton mechanics

The Lagrangian is given by: $\mathcal{L}=\sum T\left(\dot{q}_{i}\right)-V\left(q_{i}\right)$. The Hamiltonian is given by: $H=\sum \dot{q}_{i} p_{i}-\mathcal{L}$. In 2 dimensions holds: $\mathcal{L}=T-U=\frac{1}{2} m\left(\dot{r}^{2}+r^{2} \dot{\phi}^{2}\right)-U(r, \phi)$.
If the used coordinates are canonical the Hamilton equations are the equations of motion for the system:

$$
\frac{d q_{i}}{d t}=\frac{\partial H}{\partial p_{i}} ; \quad \frac{d p_{i}}{d t}=-\frac{\partial H}{\partial q_{i}}
$$

Coordinates are canonical if the following holds: $\left\{q_{i}, q_{j}\right\}=0,\left\{p_{i}, p_{j}\right\}=0,\left\{q_{i}, p_{j}\right\}=\delta_{i j}$ where $\{$,$\} is the$ Poisson bracket:

$$
\{A, B\}=\sum_{i}\left[\frac{\partial A}{\partial q_{i}} \frac{\partial B}{\partial p_{i}}-\frac{\partial A}{\partial p_{i}} \frac{\partial B}{\partial q_{i}}\right]
$$

The Hamiltonian of a Harmonic oscillator is given by $H(x, p)=p^{2} / 2 m+\frac{1}{2} m \omega^{2} x^{2}$. With new coordinates $(\theta, I)$, obtained by the canonical transformation $x=\sqrt{2 I / m \omega} \cos (\theta)$ and $p=-\sqrt{2 \operatorname{Im\omega }} \sin (\theta)$, with inverse $\theta=\arctan (-p / m \omega x)$ and $I=p^{2} / 2 m \omega+\frac{1}{2} m \omega x^{2}$ it follows: $H(\theta, I)=\omega I$.

The Hamiltonian of a charged particle with charge $q$ in an external electromagnetic field is given by:

$$
H=\frac{1}{2 m}(\vec{p}-q \vec{A})^{2}+q V
$$

This Hamiltonian can be derived from the Hamiltonian of a free particle $H=p^{2} / 2 m$ with the transformations $\vec{p} \rightarrow \vec{p}-q \vec{A}$ and $H \rightarrow H-q V$. This is elegant from a relativistic point of view: this is equivalent to the transformation of the momentum 4-vector $p^{\alpha} \rightarrow p^{\alpha}-q A^{\alpha}$. A gauge transformation on the potentials $A^{\alpha}$ corresponds with a canonical transformation, which make the Hamilton equations the equations of motion for the system.

### 1.7.3 Motion around an equilibrium, linearization

For natural systems around equilibrium the following equations are valid:

$$
\left(\frac{\partial V}{\partial q_{i}}\right)_{0}=0 ; \quad V(q)=V(0)+V_{i k} q_{i} q_{k} \text { with } V_{i k}=\left(\frac{\partial^{2} V}{\partial q_{i} \partial q_{k}}\right)_{0}
$$

With $T=\frac{1}{2}\left(M_{i k} \dot{q}_{i} \dot{q}_{k}\right)$ one receives the set of equations $M \ddot{q}+V q=0$. If $q_{i}(t)=a_{i} \exp (i \omega t)$ is substituted, this set of equations has solutions if $\operatorname{det}\left(V-\omega^{2} M\right)=0$. This leads to the eigenfrequencies of the problem: $\omega_{k}^{2}=\frac{a_{k}^{\mathrm{T}} V a_{k}}{a_{k}^{\mathrm{T}} M a_{k}}$. If the equilibrium is stable holds: $\forall k$ that $\omega_{k}^{2}>0$. The general solution is a superposition if eigenvibrations.

### 1.7.4 Phase space, Liouville's equation

In phase space holds:

$$
\nabla=\left(\sum_{i} \frac{\partial}{\partial q_{i}}, \sum_{i} \frac{\partial}{\partial p_{i}}\right) \text { so } \nabla \cdot \vec{v}=\sum_{i}\left(\frac{\partial}{\partial q_{i}} \frac{\partial H}{\partial p_{i}}-\frac{\partial}{\partial p_{i}} \frac{\partial H}{\partial q_{i}}\right)
$$

If the equation of continuity, $\partial_{t} \varrho+\nabla \cdot(\varrho \vec{v})=0$ holds, this can be written as:

$$
\{\varrho, H\}+\frac{\partial \varrho}{\partial t}=0
$$

For an arbitrary quantity $A$ holds:

$$
\frac{d A}{d t}=\{A, H\}+\frac{\partial A}{\partial t}
$$

Liouville's theorem can than be written as:

$$
\frac{d \varrho}{d t}=0 ; \quad \text { or: } \int p d q=\text { constant }
$$

### 1.7.5 Generating functions

Starting with the coordinate transformation:

$$
\left\{\begin{array}{l}
Q_{i}=Q_{i}\left(q_{i}, p_{i}, t\right) \\
P_{i}=P_{i}\left(q_{i}, p_{i}, t\right)
\end{array}\right.
$$

one can derive the following Hamilton equations with the new Hamiltonian $K$ :

$$
\frac{d Q_{i}}{d t}=\frac{\partial K}{\partial P_{i}} ; \quad \frac{d P_{i}}{d t}=-\frac{\partial K}{\partial Q_{i}}
$$

Now, a distinction between 4 cases can be made:

1. If $p_{i} \dot{q}_{i}-H=P_{i} Q_{i}-K\left(P_{i}, Q_{i}, t\right)-\frac{d F_{1}\left(q_{i}, Q_{i}, t\right)}{d t}$, the coordinates follow from:

$$
p_{i}=\frac{\partial F_{1}}{\partial q_{i}} ; \quad P_{i}=-\frac{\partial F_{1}}{\partial Q_{i}} ; \quad K=H+\frac{\partial F_{1}}{\partial t}
$$

2. If $p_{i} \dot{q}_{i}-H=-\dot{P}_{i} Q_{i}-K\left(P_{i}, Q_{i}, t\right)+\frac{d F_{2}\left(q_{i}, P_{i}, t\right)}{d t}$, the coordinates follow from:

$$
p_{i}=\frac{\partial F_{2}}{\partial q_{i}} ; \quad Q_{i}=\frac{\partial F_{2}}{\partial P_{i}} ; \quad K=H+\frac{\partial F_{2}}{\partial t}
$$

3. If $-\dot{p}_{i} q_{i}-H=P_{i} \dot{Q}_{i}-K\left(P_{i}, Q_{i}, t\right)+\frac{d F_{3}\left(p_{i}, Q_{i}, t\right)}{d t}$, the coordinates follow from:

$$
q_{i}=-\frac{\partial F_{3}}{\partial p_{i}} ; \quad P_{i}=-\frac{\partial F_{3}}{\partial Q_{i}} ; \quad K=H+\frac{\partial F_{3}}{\partial t}
$$

4. If $-\dot{p}_{i} q_{i}-H=-P_{i} Q_{i}-K\left(P_{i}, Q_{i}, t\right)+\frac{d F_{4}\left(p_{i}, P_{i}, t\right)}{d t}$, the coordinates follow from:

$$
q_{i}=-\frac{\partial F_{4}}{\partial p_{i}} ; \quad Q_{i}=\frac{\partial F_{4}}{\partial p_{i}} ; \quad K=H+\frac{\partial F_{4}}{\partial t}
$$

The functions $F_{1}, F_{2}, F_{3}$ and $F_{4}$ are called generating functions.

## Chapter 2

## Electricity \& Magnetism

### 2.1 The Maxwell equations

The classical electromagnetic field can be described by the Maxwell equations. Those can be written both as differential and integral equations:

$$
\begin{array}{ll}
\oiint(\vec{D} \cdot \vec{n}) d^{2} A=Q_{\text {free,included }} & \nabla \cdot \vec{D}=\rho_{\text {free }} \\
\oiint(\vec{B} \cdot \vec{n}) d^{2} A=0 & \nabla \cdot \vec{B}=0 \\
\oint \vec{E} \cdot d \vec{s}=-\frac{d \Phi}{d t} & \nabla \times \vec{E}=-\frac{\partial \vec{B}}{\partial t} \\
\oint \vec{H} \cdot d \vec{s}=I_{\text {free,included }}+\frac{d \Psi}{d t} & \nabla \times \vec{H}=\vec{J}_{\text {free }}+\frac{\partial \vec{D}}{\partial t}
\end{array}
$$

For the fluxes holds: $\Psi=\iint(\vec{D} \cdot \vec{n}) d^{2} A, \Phi=\iint(\vec{B} \cdot \vec{n}) d^{2} A$.
The electric displacement $\vec{D}$, polarization $\vec{P}$ and electric field strength $\vec{E}$ depend on each other according to:
$\vec{D}=\varepsilon_{0} \vec{E}+\vec{P}=\varepsilon_{0} \varepsilon_{\mathrm{r}} \vec{E}, \quad \vec{P}=\sum \vec{p}_{0} / \mathrm{Vol}, \varepsilon_{\mathrm{r}}=1+\chi_{\mathrm{e}}$, with $\chi_{\mathrm{e}}=\frac{n p_{0}^{2}}{3 \varepsilon_{0} k T}$
The magnetic field strength $\vec{H}$, the magnetization $\vec{M}$ and the magnetic flux density $\vec{B}$ depend on each other according to:
$\vec{B}=\mu_{0}(\vec{H}+\vec{M})=\mu_{0} \mu_{\mathrm{r}} \vec{H}, \vec{M}=\sum \vec{m} / \operatorname{Vol}, \mu_{\mathrm{r}}=1+\chi_{\mathrm{m}}$, with $\chi_{\mathrm{m}}=\frac{\mu_{0} n m_{0}^{2}}{3 k T}$

### 2.2 Force and potential

The force and the electric field between 2 point charges are given by:

$$
\vec{F}_{12}=\frac{Q_{1} Q_{2}}{4 \pi \varepsilon_{0} \varepsilon_{\mathrm{r}} r^{2}} \vec{e}_{r} ; \quad \vec{E}=\frac{\vec{F}}{Q}
$$

The Lorentzforce is the force which is felt by a charged particle that moves through a magnetic field. The origin of this force is a relativistic transformation of the Coulomb force: $\vec{F}_{\mathrm{L}}=Q(\vec{v} \times \vec{B})=l(\vec{I} \times \vec{B})$.
The magnetic field in point $P$ which results from an electric current is given by the law of Biot-Savart, also known als the law of Laplace. In here, $d \vec{l} \| \vec{I}$ and $\vec{r}$ points from $d \vec{l}$ to $P$ :

$$
d \vec{B}_{P}=\frac{\mu_{0} I}{4 \pi r^{2}} d \vec{l} \times \vec{e}_{r}
$$

If the current is time-dependent one has to take retardation into account: the substitution $I(t) \rightarrow I(t-r / c)$ has to be applied.
The potentials are given by: $V_{12}=-\int_{1}^{2} \vec{E} \cdot d \vec{s}$ and $\vec{A}=\frac{1}{2} \vec{B} \times \vec{r}$.

Here, the freedom remains to apply a gauge transformation. The fields can be derived from the potentials as follows:

$$
\vec{E}=-\nabla V-\frac{\partial \vec{A}}{\partial t}, \quad \vec{B}=\nabla \times \vec{A}
$$

Further holds the relation: $c^{2} \vec{B}=\vec{v} \times \vec{E}$.

### 2.3 Gauge transformations

The potentials of the electromagnetic fields transform as follows when a gauge transformation is applied:

$$
\left\{\begin{array}{l}
\overrightarrow{A^{\prime}}=\vec{A}-\nabla f \\
V^{\prime}=V+\frac{\partial f}{\partial t}
\end{array}\right.
$$

so the fields $\vec{E}$ and $\vec{B}$ do not change. This results in a canonical transformation of the Hamiltonian. Further, the freedom remains to apply a limiting condition. Two common choices are:

1. Lorentz-gauge: $\nabla \cdot \vec{A}+\frac{1}{c^{2}} \frac{\partial V}{\partial t}=0$. This separates the differential equations for $\vec{A}$ and $V: \square V=-\frac{\rho}{\varepsilon_{0}}$, $\square \vec{A}=-\mu_{0} \vec{J}$.
2. Coulomb gauge: $\nabla \cdot \vec{A}=0$. If $\rho=0$ and $\vec{J}=0$ holds $V=0$ and follows $\vec{A}$ from $\square \vec{A}=0$.

### 2.4 Energy of the electromagnetic field

The energy density of the electromagnetic field is:

$$
\frac{d W}{d \mathrm{Vol}}=w=\int H d B+\int E d D
$$

The energy density can be expressed in the potentials and currents as follows:

$$
w_{\mathrm{mag}}=\frac{1}{2} \int \vec{J} \cdot \vec{A} d^{3} x, w_{\mathrm{el}}=\frac{1}{2} \int \rho V d^{3} x
$$

### 2.5 Electromagnetic waves

### 2.5.1 Electromagnetic waves in vacuum

The wave equation $\square \Psi(\vec{r}, t)=-f(\vec{r}, t)$ has the general solution, with $c=\left(\varepsilon_{0} \mu_{0}\right)^{-1 / 2}$ :

$$
\Psi(\vec{r}, t)=\int \frac{f\left(\vec{r}, t-\left|\vec{r}-\vec{r}^{\prime}\right| / c\right)}{4 \pi\left|\vec{r}-\vec{r}^{\prime}\right|} d^{3} r^{\prime}
$$

If this is written as: $\vec{J}(\vec{r}, t)=\vec{J}(\vec{r}) \exp (-i \omega t)$ and $\vec{A}(\vec{r}, t)=\vec{A}(\vec{r}) \exp (-i \omega t)$ with:

$$
\vec{A}(\vec{r})=\frac{\mu}{4 \pi} \int \vec{J}\left(\vec{r}^{\prime}\right) \frac{\exp \left(i k\left|\vec{r}-\vec{r}^{\prime}\right|\right)}{\left|\vec{r}-\vec{r}^{\prime}\right|} d^{3} \vec{r}^{\prime} \quad, \quad V(\vec{r})=\frac{1}{4 \pi \varepsilon} \int \rho\left(\vec{r}^{\prime}\right) \frac{\exp \left(i k\left|\vec{r}-\vec{r}^{\prime}\right|\right)}{\left|\vec{r}-\vec{r}^{\prime}\right|} d^{3} \vec{r}^{\prime}
$$

A derivation via multipole expansion will show that for the radiated energy holds, if $d, \lambda \gg r$ :

$$
\frac{d P}{d \Omega}=\frac{k^{2}}{32 \pi^{2} \varepsilon_{0} c}\left|\int J_{\perp}\left(\vec{r}^{\prime}\right) \mathrm{e}^{i \vec{k} \cdot \vec{r}} d^{3} r^{\prime}\right|^{2}
$$

The energy density of the electromagnetic wave of a vibrating dipole at a large distance is:

$$
w=\varepsilon_{0} E^{2}=\frac{p_{0}^{2} \sin ^{2}(\theta) \omega^{4}}{16 \pi^{2} \varepsilon_{0} r^{2} c^{4}} \sin ^{2}(k r-\omega t), \quad\langle w\rangle_{t}=\frac{p_{0}^{2} \sin ^{2}(\theta) \omega^{4}}{32 \pi^{2} \varepsilon_{0} r^{2} c^{4}}, \quad P=\frac{c k^{4}|\vec{p}|^{2}}{12 \pi \varepsilon_{0}}
$$

The radiated energy can be derived from the Poynting vector $\vec{S}: \vec{S}=\vec{E} \times \vec{H}=c W \vec{e}_{v}$. The irradiance is the time-averaged of the Poynting vector: $I=\langle | \vec{S}| \rangle_{t}$. The radiation pressure $p_{\mathrm{s}}$ is given by $p_{\mathrm{s}}=(1+R)|\vec{S}| / c$, where $R$ is the coefficient of reflection.

### 2.5.2 Electromagnetic waves in matter

The wave equations in matter, with $c_{\text {mat }}=(\varepsilon \mu)^{-1 / 2}$ the lightspeed in matter, are:

$$
\left(\nabla^{2}-\varepsilon \mu \frac{\partial^{2}}{\partial t^{2}}-\frac{\mu}{\rho} \frac{\partial}{\partial t}\right) \vec{E}=0, \quad\left(\nabla^{2}-\varepsilon \mu \frac{\partial^{2}}{\partial t^{2}}-\frac{\mu}{\rho} \frac{\partial}{\partial t}\right) \vec{B}=0
$$

give, after substitution of monochromatic plane waves: $\vec{E}=E \exp (i(\vec{k} \cdot \vec{r}-\omega t))$ and $\vec{B}=B \exp (i(\vec{k} \cdot \vec{r}-\omega t))$ the dispersion relation:

$$
k^{2}=\varepsilon \mu \omega^{2}+\frac{i \mu \omega}{\rho}
$$

The first term arises from the displacement current, the second from the conductance current. If $k$ is written in the form $k:=k^{\prime}+i k^{\prime \prime}$ it follows that:

$$
k^{\prime}=\omega \sqrt{\frac{1}{2} \varepsilon \mu} \sqrt{1+\sqrt{1+\frac{1}{(\rho \varepsilon \omega)^{2}}}} \text { and } k^{\prime \prime}=\omega \sqrt{\frac{1}{2} \varepsilon \mu} \sqrt{-1+\sqrt{1+\frac{1}{(\rho \varepsilon \omega)^{2}}}}
$$

This results in a damped wave: $\vec{E}=E \exp \left(-k^{\prime \prime} \vec{n} \cdot \vec{r}\right) \exp \left(i\left(k^{\prime} \vec{n} \cdot \vec{r}-\omega t\right)\right)$. If the material is a good conductor, the wave vanishes after approximately one wavelength, $k=(1+i) \sqrt{\frac{\mu \omega}{2 \rho}}$.

### 2.6 Multipoles

Because $\frac{1}{\left|\vec{r}-\vec{r}^{\prime}\right|}=\frac{1}{r} \sum_{0}^{\infty}\left(\frac{r^{\prime}}{r}\right)^{l} P_{l}(\cos \theta)$ the potential can be written as: $V=\frac{Q}{4 \pi \varepsilon} \sum_{n} \frac{k_{n}}{r^{n}}$
For the lowest-order terms this results in:

- Monopole: $l=0, k_{0}=\int \rho d V$
- Dipole: $l=1, k_{1}=\int r \cos (\theta) \rho d V$
- Quadrupole: $l=2, k_{2}=\frac{1}{2} \sum_{i}\left(3 z_{i}^{2}-r_{i}^{2}\right)$

1. The electric dipole: dipole moment: $\vec{p}=Q l \vec{e}$, where $\vec{e}$ goes from $\oplus$ to $\ominus$, and $\vec{F}=(\vec{p} \cdot \nabla) \vec{E}_{\text {ext }}$, and $W=-\vec{p} \cdot \vec{E}_{\text {out }}$.
Electric field: $\vec{E} \approx \frac{Q}{4 \pi \varepsilon r^{3}}\left(\frac{3 \vec{p} \cdot \vec{r}}{r^{2}}-\vec{p}\right)$. The torque is: $\vec{\tau}=\vec{p} \times \vec{E}_{\text {out }}$
2. The magnetic dipole: dipole moment: if $r \gg \sqrt{A}: \vec{\mu}=\vec{I} \times\left(A \vec{e}_{\perp}\right), \vec{F}=(\vec{\mu} \cdot \nabla) \vec{B}_{\text {out }}$
$|\mu|=\frac{m v_{\perp}^{2}}{2 B}, W=-\vec{\mu} \times \vec{B}_{\text {out }}$
Magnetic field: $\vec{B}=\frac{-\mu}{4 \pi r^{3}}\left(\frac{3 \mu \cdot \vec{r}}{r^{2}}-\vec{\mu}\right)$. The moment is: $\vec{\tau}=\vec{\mu} \times \vec{B}_{\text {out }}$

### 2.7 Electric currents

The continuity equation for charge is: $\frac{\partial \rho}{\partial t}+\nabla \cdot \vec{J}=0$. The electric current is given by:

$$
I=\frac{d Q}{d t}=\iint(\vec{J} \cdot \vec{n}) d^{2} A
$$

For most conductors holds: $\vec{J}=\vec{E} / \rho$, where $\rho$ is the resistivity.

If the flux enclosed by a conductor changes this results in an induced voltage $V_{\mathrm{ind}}=-N \frac{d \Phi}{d t}$. If the current flowing through a conductor changes, this results in a self-inductance which opposes the original change: $V_{\text {selfind }}=-L \frac{d I}{d t}$. If a conductor encloses a flux $\Phi$ holds: $\Phi=L I$.
The magnetic induction within a coil is approximated by: $B=\frac{\mu N I}{\sqrt{l^{2}+4 R^{2}}}$ where $l$ is the length, $R$ the radius and $N$ the number of coils. The energy contained within a coil is given by $W=\frac{1}{2} L I^{2}$ and $L=\mu N^{2} A / l$.
The capacity is defined by: $C=Q / V$. For a capacitor holds: $C=\varepsilon_{0} \varepsilon_{\mathrm{r}} A / d$ where $d$ is the distance between the plates and $A$ the surface of one plate. The electric field strength between the plates is $E=\sigma / \varepsilon_{0}=Q / \varepsilon_{0} A$ where $\sigma$ is the surface charge. The accumulated energy is given by $W=\frac{1}{2} C V^{2}$. The current through a capacity is given by $I=-C \frac{d V}{d t}$.
For most PTC resistors holds approximately: $R=R_{0}(1+\alpha T)$, where $R_{0}=\rho l / A$. For a NTC holds: $R(T)=C \exp (-B / T)$ where $B$ and $C$ depend only on the material.
If a current flows through two different, connecting conductors $x$ and $y$, the contact area will heat up or cool down, depending on the direction of the current: the Peltier effect. The generated or removed heat is given by: $W=\Pi_{x y} I t$. This effect can be amplified with semiconductors.
The thermic voltage between 2 metals is given by: $V=\gamma\left(T-T_{0}\right)$. For a Cu-Konstantane connection holds: $\gamma \approx 0.2-0.7 \mathrm{mV} / \mathrm{K}$.

In an electrical net with only stationary currents, Kirchhoff's equations apply: for a knot holds: $\sum I_{n}=0$, along a closed path holds: $\sum V_{n}=\sum I_{n} R_{n}=0$.

### 2.8 Depolarizing field

If a dielectric material is placed in an electric or magnetic field, the field strength within and outside the material will change because the material will be polarized or magnetized. If the medium has an ellipsoidal shape and one of the principal axes is parallel with the external field $\vec{E}_{0}$ or $\vec{B}_{0}$ then the depolarizing is field homogeneous.

$$
\begin{gathered}
\vec{E}_{\mathrm{dep}}=\vec{E}_{\mathrm{mat}}-\vec{E}_{0}=-\frac{\mathcal{N} \vec{P}}{\varepsilon_{0}} \\
\vec{H}_{\mathrm{dep}}=\vec{H}_{\mathrm{mat}}-\vec{H}_{0}=-\mathcal{N} \vec{M}
\end{gathered}
$$

$\mathcal{N}$ is a constant depending only on the shape of the object placed in the field, with $0 \leq \mathcal{N} \leq 1$. For a few limiting cases of an ellipsoid holds: a thin plane: $\mathcal{N}=1$, a long, thin bar: $\mathcal{N}=0$, a sphere: $\mathcal{N}=\frac{1}{3}$.

### 2.9 Mixtures of materials

The average electric displacement in a material which is inhomogenious on a mesoscopic scale is given by: $\langle D\rangle=\langle\varepsilon E\rangle=\varepsilon^{*}\langle E\rangle$ where $\varepsilon^{*}=\varepsilon_{1}\left(1-\frac{\phi_{2}(1-x)}{\Phi\left(\varepsilon^{*} / \varepsilon_{2}\right)}\right)^{-1}$ where $x=\varepsilon_{1} / \varepsilon_{2}$. For a sphere holds: $\Phi=$ $\frac{1}{3}+\frac{2}{3} x$. Further holds:

$$
\left(\sum_{i} \frac{\phi_{i}}{\varepsilon_{i}}\right)^{-1} \leq \varepsilon^{*} \leq \sum_{i} \phi_{i} \varepsilon_{i}
$$

## Chapter 3

## Relativity

### 3.1 Special relativity

### 3.1.1 The Lorentz transformation

The Lorentz transformation $\left(\vec{x}^{\prime}, t^{\prime}\right)=\left(\vec{x}^{\prime}(\vec{x}, t), t^{\prime}(\vec{x}, t)\right)$ leaves the wave equation invariant if $c$ is invariant:

$$
\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}+\frac{\partial^{2}}{\partial z^{2}}-\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}=\frac{\partial^{2}}{\partial x^{\prime 2}}+\frac{\partial^{2}}{\partial y^{\prime 2}}+\frac{\partial^{2}}{\partial z^{\prime 2}}-\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{\prime 2}}
$$

This transformation can also be found when $d s^{2}=d s^{\prime 2}$ is demanded. The general form of the Lorentz transformation is given by:

$$
\vec{x}^{\prime}=\vec{x}+\frac{(\gamma-1)(\vec{x} \cdot \vec{v}) \vec{v}}{|v|^{2}}-\gamma \vec{v} t, \quad t^{\prime}=\gamma\left(t-\frac{\vec{x} \cdot \vec{v}}{c^{2}}\right)
$$

where

$$
\gamma=\frac{1}{\sqrt{1-\frac{v^{2}}{c^{2}}}}
$$

The velocity difference $\vec{v}^{\prime}$ between two observers transforms according to:

$$
\vec{v}^{\prime}=\left(\gamma\left(1-\frac{\vec{v}_{1} \cdot \vec{v}_{2}}{c^{2}}\right)\right)^{-1}\left(\vec{v}_{2}+(\gamma-1) \frac{\vec{v}_{1} \cdot \vec{v}_{2}}{v_{1}^{2}} \vec{v}_{1}-\gamma \vec{v}_{1}\right)
$$

If the velocity is parallel to the $x$-axis, this becomes $y^{\prime}=y, z^{\prime}=z$ and:

$$
\begin{aligned}
& x^{\prime}=\gamma(x-v t), \quad x=\gamma\left(x^{\prime}+v t^{\prime}\right) \\
& t^{\prime}=\gamma\left(t-\frac{x v}{c^{2}}\right), \quad t=\gamma\left(t^{\prime}+\frac{x^{\prime} v}{c^{2}}\right), \quad v^{\prime}=\frac{v_{2}-v_{1}}{1-\frac{v_{1} v_{2}}{c^{2}}}
\end{aligned}
$$

If $\vec{v}=v \vec{e}_{x}$ holds:

$$
p_{x}^{\prime}=\gamma\left(p_{x}-\frac{\beta W}{c}\right), \quad W^{\prime}=\gamma\left(W-v p_{x}\right)
$$

With $\beta=v / c$ the electric field of a moving charge is given by:

$$
\vec{E}=\frac{Q}{4 \pi \varepsilon_{0} r^{2}} \frac{\left(1-\beta^{2}\right) \vec{e}_{r}}{\left(1-\beta^{2} \sin ^{2}(\theta)\right)^{3 / 2}}
$$

The electromagnetic field transforms according to:

$$
\vec{E}^{\prime}=\gamma(\vec{E}+\vec{v} \times \vec{B}) \quad, \quad \vec{B}^{\prime}=\gamma\left(\vec{B}-\frac{\vec{v} \times \vec{E}}{c^{2}}\right)
$$

Length, mass and time transform according to: $\Delta t_{\mathrm{r}}=\gamma \Delta t_{0}, m_{\mathrm{r}}=\gamma m_{0}, l_{\mathrm{r}}=l_{0} / \gamma$, with ${ }_{0}$ the quantities in a co-moving reference frame and ${ }_{\mathrm{r}}$ the quantities in a frame moving with velocity $v$ w.r.t. it. The proper time $\tau$ is defined as: $d \tau^{2}=d s^{2} / c^{2}$, so $\Delta \tau=\Delta t / \gamma$. For energy and momentum holds: $W=m_{\mathrm{r}} c^{2}=\gamma W_{0}$,
$W^{2}=m_{0}^{2} c^{4}+p^{2} c^{2} . p=m_{\mathrm{r}} v=\gamma m_{0} v=W v / c^{2}$, and $p c=W \beta$ where $\beta=v / c$. The force is defined by $\vec{F}=d \vec{p} / d t$.
4-vectors have the property that their modulus is independent of the observer: their components can change after a coordinate transformation but not their modulus. The difference of two 4 -vectors transforms also as a 4 -vector. The 4 -vector for the velocity is given by $U^{\alpha}=\frac{d x^{\alpha}}{d \tau}$. The relation with the "common" velocity $u^{i}:=d x^{i} / d t$ is: $U^{\alpha}=\left(\gamma u^{i}, i c \gamma\right)$. For particles with nonzero restmass holds: $U^{\alpha} U_{\alpha}=-c^{2}$, for particles with zero restmass (so with $v=c$ ) holds: $U^{\alpha} U_{\alpha}=0$. The 4-vector for energy and momentum is given by: $p^{\alpha}=m_{0} U^{\alpha}=\left(\gamma p^{i}, i W / c\right)$. So: $p_{\alpha} p^{\alpha}=-m_{0}^{2} c^{2}=p^{2}-W^{2} / c^{2}$.

### 3.1.2 Red and blue shift

There are three causes of red and blue shifts:

1. Motion: with $\vec{e}_{v} \cdot \vec{e}_{r}=\cos (\varphi)$ follows: $\frac{f^{\prime}}{f}=\gamma\left(1-\frac{v \cos (\varphi)}{c}\right)$.

This can give both red- and blueshift, also $\perp$ to the direction of motion.
2. Gravitational redshift: $\frac{\Delta f}{f}=\frac{\kappa M}{r c^{2}}$.
3. Redshift because the universe expands, resulting in e.g. the cosmic background radiation:

$$
\frac{\lambda_{0}}{\lambda_{1}}=\frac{R_{0}}{R_{1}} .
$$

### 3.1.3 The stress-energy tensor and the field tensor

The stress-energy tensor is given by:

$$
T_{\mu \nu}=\left(\varrho c^{2}+p\right) u_{\mu} u_{\nu}+p g_{\mu \nu}+\frac{1}{c^{2}}\left(F_{\mu \alpha} F_{\nu}^{\alpha}+\frac{1}{4} g_{\mu \nu} F^{\alpha \beta} F_{\alpha \beta}\right)
$$

The conservation laws can than be written as: $\nabla_{\nu} T^{\mu \nu}=0$. The electromagnetic field tensor is given by:

$$
F_{\alpha \beta}=\frac{\partial A_{\beta}}{\partial x^{\alpha}}-\frac{\partial A_{\alpha}}{\partial x^{\beta}}
$$

with $A_{\mu}:=(\vec{A}, i V / c)$ and $J_{\mu}:=(\vec{J}, i c \rho)$. The Maxwell equations can than be written as:

$$
\partial_{\nu} F^{\mu \nu}=\mu_{0} J^{\mu}, \partial_{\lambda} F_{\mu \nu}+\partial_{\mu} F_{\nu \lambda}+\partial_{\nu} F_{\lambda \mu}=0
$$

The equations of motion for a charged particle in an EM field become with the field tensor:

$$
\frac{d p_{\alpha}}{d \tau}=q F_{\alpha \beta} u^{\beta}
$$

### 3.2 General relativity

### 3.2.1 Riemannian geometry, the Einstein tensor

The basic principles of general relativity are:

1. The geodesic postulate: free falling particles move along geodesics of space-time with the proper time $\tau$ or arc length $s$ as parameter. For particles with zero rest mass (photons), the use of a free parameter is required because for them holds $d s=0$. From $\delta \int d s=0$ the equations of motion can be derived:

$$
\frac{d^{2} x^{\alpha}}{d s^{2}}+\Gamma_{\beta \gamma}^{\alpha} \frac{d x^{\beta}}{d s} \frac{d x^{\gamma}}{d s}=0
$$

2. The principle of equivalence: inertial mass $\equiv$ gravitational mass $\Rightarrow$ gravitation is equivalent with a curved space-time were particles move along geodesics.
3. By a proper choice of the coordinate system it is possible to make the metric locally flat in each point $x_{i}: g_{\alpha \beta}\left(x_{i}\right)=\eta_{\alpha \beta}:=\operatorname{diag}(-1,1,1,1)$.

The Riemann tensor is defined as: $R_{\nu \alpha \beta}^{\mu} T^{\nu}:=\nabla_{\alpha} \nabla_{\beta} T^{\mu}-\nabla_{\beta} \nabla_{\alpha} T^{\mu}$, where the covariant derivative is given by $\nabla_{j} a^{i}=\partial_{j} a^{i}+\Gamma_{j k}^{i} a^{k}$ and $\nabla_{j} a_{i}=\partial_{j} a_{i}-\Gamma_{i j}^{k} a_{k}$. Here,

$$
\Gamma_{j k}^{i}=\frac{g^{i l}}{2}\left(\frac{\partial g_{l j}}{\partial x^{k}}+\frac{\partial g_{l k}}{\partial x^{j}}-\frac{\partial g_{j} k}{\partial x^{l}}\right), \text { for Euclidean spaces this reduces to: } \Gamma_{j k}^{i}=\frac{\partial^{2} \bar{x}^{l}}{\partial x^{j} \partial x^{k}} \frac{\partial x^{i}}{\partial \bar{x}^{l}},
$$

are the Christoffel symbols. For a second-order tensor holds: $\left[\nabla_{\alpha}, \nabla_{\beta}\right] T_{\nu}^{\mu}=R_{\sigma \alpha \beta}^{\mu} T_{\nu}^{\sigma}+R_{\nu \alpha \beta}^{\sigma} T_{\sigma}^{\mu}, \nabla_{k} a_{j}^{i}=$ $\partial_{k} a_{j}^{i}-\Gamma_{k j}^{l} a_{l}^{i}+\Gamma_{k l}^{i} a_{j}^{l}, \nabla_{k} a_{i j}=\partial_{k} a_{i j}-\Gamma_{k i}^{l} a_{l j}-\Gamma_{k j}^{l} a_{j l}$ and $\nabla_{k} a^{i j}=\partial_{k} a^{i j}+\Gamma_{k l}^{i} a^{l j}+\Gamma_{k l}^{j} a^{i l}$. The following holds: $R_{\beta \mu \nu}^{\alpha}=\partial_{\mu} \Gamma_{\beta \nu}^{\alpha}-\partial_{\nu} \Gamma_{\beta \mu}^{\alpha}+\Gamma_{\sigma \mu}^{\alpha} \Gamma_{\beta \nu}^{\sigma}-\Gamma_{\sigma \nu}^{\alpha} \Gamma_{\beta \mu}^{\sigma}$.

The Ricci tensor is a contraction of the Riemann tensor: $R_{\alpha \beta}:=R_{\alpha \mu \beta}^{\mu}$, which is symmetric: $R_{\alpha \beta}=R_{\beta \alpha}$. The Bianchi identities are: $\nabla_{\lambda} R_{\alpha \beta \mu \nu}+\nabla_{\nu} R_{\alpha \beta \lambda \mu}+\nabla_{\mu} R_{\alpha \beta \nu \lambda}=0$.
The Einstein tensor is given by: $G^{\alpha \beta}:=R^{\alpha \beta}-\frac{1}{2} g^{\alpha \beta} R$, where $R:=R_{\alpha}^{\alpha}$ is the Ricci scalar, for which holds: $\nabla_{\beta} G_{\alpha \beta}=0$. With the variational principle $\delta \int\left(\mathcal{L}\left(g_{\mu \nu}\right)-R c^{2} / 16 \pi \kappa\right) \sqrt{|g|} d^{4} x=0$ for variations $g_{\mu \nu} \rightarrow g_{\mu \nu}+\delta g_{\mu \nu}$ the Einstein field equations can be derived:

$$
G_{\alpha \beta}=\frac{8 \pi \kappa}{c^{2}} T_{\alpha \beta}, \text { which can also be written as } \quad R_{\alpha \beta}=\frac{8 \pi \kappa}{c^{2}}\left(T_{\alpha \beta}-\frac{1}{2} g_{\alpha \beta} T_{\mu}^{\mu}\right)
$$

For empty space this is equivalent to $R_{\alpha \beta}=0$. The equation $R_{\alpha \beta \mu \nu}=0$ has as only solution a flat space.
The Einstein equations are 10 independent equations, which are of second order in $g_{\mu \nu}$. From this, the Laplace equation from Newtonian gravitation can be derived by stating: $g_{\mu \nu}=\eta_{\mu \nu}+h_{\mu \nu}$, where $|h| \ll 1$. In the stationary case, this results in $\nabla^{2} h_{00}=8 \pi \kappa \varrho / c^{2}$.

The most general form of the field equations is: $R_{\alpha \beta}-\frac{1}{2} g_{\alpha \beta} R+\Lambda g_{\alpha \beta}=\frac{8 \pi \kappa}{c^{2}} T_{\alpha \beta}$
where $\Lambda$ is the cosmological constant. This constant plays a role in inflatory models of the universe.

### 3.2.2 The line element

The metric tensor in an Euclidean space is given by: $g_{i j}=\sum_{k} \frac{\partial \bar{x}^{k}}{\partial x^{i}} \frac{\partial \bar{x}^{k}}{\partial x^{j}}$.
In general holds: $d s^{2}=g_{\mu \nu} d x^{\mu} d x^{\nu}$. In special relativity this becomes $d s^{2}=-c^{2} d t^{2}+d x^{2}+d y^{2}+d z^{2}$. This metric, $\eta_{\mu \nu}:=\operatorname{diag}(-1,1,1,1)$, is called the Minkowski metric.
The external Schwarzschild metric applies in vacuum outside a spherical mass distribution, and is given by:

$$
d s^{2}=\left(-1+\frac{2 m}{r}\right) c^{2} d t^{2}+\left(1-\frac{2 m}{r}\right)^{-1} d r^{2}+r^{2} d \Omega^{2}
$$

Here, $m:=M \kappa / c^{2}$ is the geometrical mass of an object with mass $M$, and $d \Omega^{2}=d \theta^{2}+\sin ^{2} \theta d \varphi^{2}$. This metric is singular for $r=2 m=2 \kappa M / c^{2}$. If an object is smaller than its event horizon $2 m$, that implies that its escape velocity is $>c$, it is called a black hole. The Newtonian limit of this metric is given by:

$$
d s^{2}=-(1+2 V) c^{2} d t^{2}+(1-2 V)\left(d x^{2}+d y^{2}+d z^{2}\right)
$$

where $V=-\kappa M / r$ is the Newtonian gravitation potential. In general relativity, the components of $g_{\mu \nu}$ are associated with the potentials and the derivatives of $g_{\mu \nu}$ with the field strength.
The Kruskal-Szekeres coordinates are used to solve certain problems with the Schwarzschild metric near $r=2 m$. They are defined by:

- $r>2 m$ :

$$
\left\{\begin{aligned}
u & =\sqrt{\frac{r}{2 m}-1} \exp \left(\frac{r}{4 m}\right) \cosh \left(\frac{t}{4 m}\right) \\
v & =\sqrt{\frac{r}{2 m}-1} \exp \left(\frac{r}{4 m}\right) \sinh \left(\frac{t}{4 m}\right)
\end{aligned}\right.
$$

- $r<2 m$ :

$$
\left\{\begin{aligned}
u & =\sqrt{1-\frac{r}{2 m}} \exp \left(\frac{r}{4 m}\right) \sinh \left(\frac{t}{4 m}\right) \\
v & =\sqrt{1-\frac{r}{2 m}} \exp \left(\frac{r}{4 m}\right) \cosh \left(\frac{t}{4 m}\right)
\end{aligned}\right.
$$

- $r=2 m$ : here, the Kruskal coordinates are singular, which is necessary to eliminate the coordinate singularity there.
The line element in these coordinates is given by:

$$
d s^{2}=-\frac{32 m^{3}}{r} \mathrm{e}^{-r / 2 m}\left(d v^{2}-d u^{2}\right)+r^{2} d \Omega^{2}
$$

The line $r=2 m$ corresponds to $u=v=0$, the limit $x^{0} \rightarrow \infty$ with $u=v$ and $x^{0} \rightarrow-\infty$ with $u=-v$. The Kruskal coordinates are only singular on the hyperbole $v^{2}-u^{2}=1$, this corresponds with $r=0$. On the line $d v= \pm d u$ holds $d \theta=d \varphi=d s=0$.

For the metric outside a rotating, charged spherical mass the Newman metric applies:

$$
\begin{aligned}
d s^{2}= & \left(1-\frac{2 m r-e^{2}}{r^{2}+a^{2} \cos ^{2} \theta}\right) c^{2} d t^{2}-\left(\frac{r^{2}+a^{2} \cos ^{2} \theta}{r^{2}-2 m r+a^{2}-e^{2}}\right) d r^{2}-\left(r^{2}+a^{2} \cos ^{2} \theta\right) d \theta^{2}- \\
& \left(r^{2}+a^{2}+\frac{\left(2 m r-e^{2}\right) a^{2} \sin ^{2} \theta}{r^{2}+a^{2} \cos ^{2} \theta}\right) \sin ^{2} \theta d \varphi^{2}+\left(\frac{2 a\left(2 m r-e^{2}\right)}{r^{2}+a^{2} \cos ^{2} \theta}\right) \sin ^{2} \theta(d \varphi)(c d t)
\end{aligned}
$$

where $m=\kappa M / c^{2}, a=L / M c$ and $e=\kappa Q / \varepsilon_{0} c^{2}$.
A rotating charged black hole has an event horizon with $R_{\mathrm{S}}=m+\sqrt{m^{2}-a^{2}-e^{2}}$.
Near rotating black holes frame dragging occurs because $g_{t \varphi} \neq 0$. For the Kerr metric ( $e=0, a \neq 0$ ) then follows that within the surface $R_{\mathrm{E}}=m+\sqrt{m^{2}-a^{2} \cos ^{2} \theta}$ (de ergosphere) no particle can be at rest.

### 3.2.3 Planetary orbits and the perihelion shift

To find a planetary orbit, the variational problem $\delta \int d s=0$ has to be solved. This is equivalent to the problem $\delta \int d s^{2}=\delta \int g_{i j} d x^{i} d x^{j}=0$. Substituting the external Schwarzschild metric yields for a planetary orbit:

$$
\frac{d u}{d \varphi}\left(\frac{d^{2} u}{d \varphi^{2}}+u\right)=\frac{d u}{d \varphi}\left(3 m u+\frac{m}{h^{2}}\right)
$$

where $u:=1 / r$ and $h=r^{2} \dot{\varphi}=$ constant. The term $3 m u$ is not present in the classical solution. This term can in the classical case also be found from a potential $V(r)=-\frac{\kappa M}{r}\left(1+\frac{h^{2}}{r^{2}}\right)$.
The orbital equation gives $r=$ constant as solution, or can, after dividing by $d u / d \varphi$, be solved with perturbation theory. In zeroth order, this results in an elliptical orbit: $u_{0}(\varphi)=A+B \cos (\varphi)$ with $A=m / h^{2}$ and $B$ an arbitrary constant. In first order, this becomes:

$$
u_{1}(\varphi)=A+B \cos (\varphi-\varepsilon \varphi)+\varepsilon\left(A+\frac{B^{2}}{2 A}-\frac{B^{2}}{6 A} \cos (2 \varphi)\right)
$$

where $\varepsilon=3 m^{2} / h^{2}$ is small. The perihelion of a planet is the point for which $r$ is minimal, or $u$ maximal. This is the case if $\cos (\varphi-\varepsilon \varphi)=0 \Rightarrow \varphi \approx 2 \pi n(1+\varepsilon)$. For the perihelion shift then follows: $\Delta \varphi=2 \pi \varepsilon=$ $6 \pi m^{2} / h^{2}$ per orbit.

### 3.2.4 The trajectory of a photon

For the trajectory of a photon (and for each particle with zero restmass) holds $d s^{2}=0$. Substituting the external Schwarzschild metric results in the following orbital equation:

$$
\frac{d u}{d \varphi}\left(\frac{d^{2} u}{d \varphi^{2}}+u-3 m u\right)=0
$$

### 3.2.5 Gravitational waves

Starting with the approximation $g_{\mu \nu}=\eta_{\mu \nu}+h_{\mu \nu}$ for weak gravitational fields and the definition $h_{\mu \nu}^{\prime}=$ $h_{\mu \nu}-\frac{1}{2} \eta_{\mu \nu} h_{\alpha}^{\alpha}$ it follows that $\square h_{\mu \nu}^{\prime}=0$ if the gauge condition $\partial h_{\mu \nu}^{\prime} / \partial x^{\nu}=0$ is satisfied. From this, it follows that the loss of energy of a mechanical system, if the occurring velocities are $\ll c$ and for wavelengths $\gg$ the size of the system, is given by:

$$
\frac{d E}{d t}=-\frac{G}{5 c^{5}} \sum_{i, j}\left(\frac{d^{3} Q_{i j}}{d t^{3}}\right)^{2}
$$

with $Q_{i j}=\int \varrho\left(x_{i} x_{j}-\frac{1}{3} \delta_{i j} r^{2}\right) d^{3} x$ the mass quadrupole moment.

### 3.2.6 Cosmology

If for the universe as a whole is assumed:

1. There exists a global time coordinate which acts as $x^{0}$ of a Gaussian coordinate system,
2. The 3-dimensional spaces are isotrope for a certain value of $x^{0}$,
3. Each point is equivalent to each other point for a fixed $x^{0}$.
then the Robertson-Walker metric can be derived for the line element:

$$
d s^{2}=-c^{2} d t^{2}+\frac{R^{2}(t)}{r_{0}^{2}\left(1-\frac{k r^{2}}{4 r_{0}^{2}}\right)}\left(d r^{2}+r^{2} d \Omega^{2}\right)
$$

For the scalefactor $R(t)$ the following equations can be derived:

$$
\frac{2 \ddot{R}}{R}+\frac{\dot{R}^{2}+k c^{2}}{R^{2}}=-\frac{8 \pi \kappa p}{c^{2}}+\Lambda \quad \text { and } \quad \frac{\dot{R}^{2}+k c^{2}}{R^{2}}=\frac{8 \pi \kappa \varrho}{3}+\frac{\Lambda}{3}
$$

where $p$ is the pressure and $\varrho$ the density of the universe. If $\Lambda=0$ can be derived for the deceleration parameter $q$ :

$$
q=-\frac{\ddot{R} R}{\dot{R}^{2}}=\frac{4 \pi \kappa \varrho}{3 H^{2}}
$$

where $H=\dot{R} / R$ is Hubble's constant. This is a measure of the velocity with which galaxies far away are moving away from each other, and has the value $\approx(75 \pm 25) \mathrm{km} \cdot \mathrm{s}^{-1} \cdot \mathrm{Mpc}^{-1}$. This gives 3 possible conditions for the universe (here, $W$ is the total amount of energy in the universe):

1. Parabolical universe: $k=0, W=0, q=\frac{1}{2}$. The expansion velocity of the universe $\rightarrow 0$ if $t \rightarrow \infty$. The hereto related critical density is $\varrho_{\mathrm{c}}=3 H^{2} / 8 \pi \kappa$.
2. Hyperbolical universe: $k=-1, W<0, q<\frac{1}{2}$. The expansion velocity of the universe remains positive forever.
3. Elliptical universe: $k=1, W>0, q>\frac{1}{2}$. The expansion velocity of the universe becomes negative after some time: the universe starts collapsing.

## Chapter 4

## Oscillations

### 4.1 Harmonic oscillations

The general form of a harmonic oscillation is: $\Psi(t)=\hat{\Psi} \mathrm{e}^{i(\omega t \pm \varphi)} \equiv \hat{\Psi} \cos (\omega t \pm \varphi)$,
where $\hat{\Psi}$ is the amplitude. A superposition of several harmonic oscillations with the same frequency results in another harmonic oscillation:

$$
\sum_{i} \hat{\Psi}_{i} \cos \left(\alpha_{i} \pm \omega t\right)=\hat{\Phi} \cos (\beta \pm \omega t)
$$

with:

$$
\tan (\beta)=\frac{\sum_{i} \hat{\Psi}_{i} \sin \left(\alpha_{i}\right)}{\sum_{i} \hat{\Psi}_{i} \cos \left(\alpha_{i}\right)} \text { and } \quad \hat{\Phi}^{2}=\sum_{i} \hat{\Psi}_{i}^{2}+2 \sum_{j>i} \sum_{i} \hat{\Psi}_{i} \hat{\Psi}_{j} \cos \left(\alpha_{i}-\alpha_{j}\right)
$$

For harmonic oscillations holds: $\int x(t) d t=\frac{x(t)}{i \omega}$ and $\frac{d^{n} x(t)}{d t^{n}}=(i \omega)^{n} x(t)$.

### 4.2 Mechanic oscillations

For a construction with a spring with constant $C$ parallel to a damping $k$ which is connected to a mass $M$, to which a periodic force $F(t)=\hat{F} \cos (\omega t)$ is applied holds the equation of motion $m \ddot{x}=F(t)-k \dot{x}-C x$. With complex amplitudes, this becomes $-m \omega^{2} x=F-C x-i k \omega x$. With $\omega_{0}^{2}=C / m$ follows:

$$
x=\frac{F}{m\left(\omega_{0}^{2}-\omega^{2}\right)+i k \omega}, \text { and for the velocity holds: } \dot{x}=\frac{F}{i \sqrt{C m} \delta+k}
$$

where $\delta=\frac{\omega}{\omega_{0}}-\frac{\omega_{0}}{\omega}$. The quantity $Z=F / \dot{x}$ is called the impedance of the system. The quality of the system is given by $Q=\frac{\sqrt{C m}}{k}$.
The frequency with minimal $|Z|$ is called velocity resonance frequency. This is equal to $\omega_{0}$. In the resonance curve $|Z| / \sqrt{C m}$ is plotted against $\omega / \omega_{0}$. The width of this curve is characterized by the points where $|Z(\omega)|=$ $\left|Z\left(\omega_{0}\right)\right| \sqrt{2}$. In these points holds: $R=X$ and $\delta= \pm Q^{-1}$, and the width is $2 \Delta \omega_{\mathrm{B}}=\omega_{0} / Q$.
The stiffness of an oscillating system is given by $F / x$. The amplitude resonance frequency $\omega_{\mathrm{A}}$ is the frequency where $i \omega Z$ is minimal. This is the case for $\omega_{\mathrm{A}}=\omega_{0} \sqrt{1-\frac{1}{2} Q^{2}}$.
The damping frequency $\omega_{\mathrm{D}}$ is a measure for the time in which an oscillating system comes to rest. It is given by $\omega_{\mathrm{D}}=\omega_{0} \sqrt{1-\frac{1}{4 Q^{2}}}$. A weak damped oscillation $\left(k^{2}<4 m C\right)$ dies out after $T_{\mathrm{D}}=2 \pi / \omega_{\mathrm{D}}$. For a critical damped oscillation $\left(k^{2}=4 m C\right)$ holds $\omega_{\mathrm{D}}=0$. A strong damped oscillation $\left(k^{2}>4 m C\right)$ drops like (if $\left.k^{2} \gg 4 m C\right) x(t) \approx x_{0} \exp (-t / \tau)$.

### 4.3 Electric oscillations

The impedance is given by: $Z=R+i X$. The phase angle is $\varphi:=\arctan (X / R)$. The impedance of a resistor is $R$, of a capacitor $1 / i \omega C$ and of a self inductor $i \omega L$. The quality of a coil is $Q=\omega L / R$. The total impedance in case several elements are positioned is given by:

1. Series connection: $V=I Z$,

$$
Z_{\mathrm{tot}}=\sum_{i} Z_{i}, \quad L_{\mathrm{tot}}=\sum_{i} L_{i}, \frac{1}{C_{\mathrm{tot}}}=\sum_{i} \frac{1}{C_{i}}, \quad Q=\frac{Z_{0}}{R}, \quad Z=R(1+i Q \delta)
$$

2. parallel connection: $V=I Z$,

$$
\frac{1}{Z_{\mathrm{tot}}}=\sum_{i} \frac{1}{Z_{i}}, \frac{1}{L_{\mathrm{tot}}}=\sum_{i} \frac{1}{L_{i}}, \quad C_{\mathrm{tot}}=\sum_{i} C_{i}, \quad Q=\frac{R}{Z_{0}}, \quad Z=\frac{R}{1+i Q \delta}
$$

Here, $Z_{0}=\sqrt{\frac{L}{C}}$ and $\omega_{0}=\frac{1}{\sqrt{L C}}$.
The power given by a source is given by $P(t)=V(t) \cdot I(t)$, so $\langle P\rangle_{t}=\hat{V}_{\text {eff }} \hat{I}_{\text {eff }} \cos (\Delta \phi)$ $=\frac{1}{2} \hat{V} \hat{I} \cos \left(\phi_{v}-\phi_{i}\right)=\frac{1}{2} \hat{I}^{2} \operatorname{Re}(Z)=\frac{1}{2} \hat{V}^{2} \operatorname{Re}(1 / Z)$, where $\cos (\Delta \phi)$ is the work factor.

### 4.4 Waves in long conductors

These cables are in use for signal transfer, e.g. coax cable. For them holds: $Z_{0}=\sqrt{\frac{d L}{d x}} \frac{d x}{d C}$.
The transmission velocity is given by $v=\sqrt{\frac{d x}{d L} \frac{d x}{d C}}$.

### 4.5 Coupled conductors and transformers

For two coils enclosing each others flux holds: if $\Phi_{12}$ is the part of the flux originating from $I_{2}$ through coil 2 which is enclosed by coil 1 , than holds $\Phi_{12}=M_{12} I_{2}, \Phi_{21}=M_{21} I_{1}$. For the coefficients of mutual induction $M_{i j}$ holds:

$$
M_{12}=M_{21}:=M=k \sqrt{L_{1} L_{2}}=\frac{N_{1} \Phi_{1}}{I_{2}}=\frac{N_{2} \Phi_{2}}{I_{1}} \sim N_{1} N_{2}
$$

where $0 \leq k \leq 1$ is the coupling factor. For a transformer is $k \approx 1$. At full load holds:

$$
\frac{V_{1}}{V_{2}}=\frac{I_{2}}{I_{1}}=-\frac{i \omega M}{i \omega L_{2}+R_{\mathrm{load}}} \approx-\sqrt{\frac{L_{1}}{L_{2}}}=-\frac{N_{1}}{N_{2}}
$$

### 4.6 Pendulums

The oscillation time $T=1 / f$, and for different types of pendulums is given by:

- Oscillating spring: $T=2 \pi \sqrt{m / C}$ if the spring force is given by $F=C \cdot \Delta l$.
- Physical pendulum: $T=2 \pi \sqrt{I / \tau}$ with $\tau$ the moment of force and $I$ the moment of inertia.
- Torsion pendulum: $T=2 \pi \sqrt{I / \kappa}$ with $\kappa=\frac{2 l m}{\pi r^{4} \Delta \varphi}$ the constant of torsion and $I$ the moment of inertia.
- Mathematical pendulum: $T=2 \pi \sqrt{l / g}$ with $g$ the acceleration of gravity and $l$ the length of the pendulum.


## Chapter 5

## Waves

### 5.1 The wave equation

The general form of the wave equation is: $\square u=0$, or:

$$
\nabla^{2} u-\frac{1}{v^{2}} \frac{\partial^{2} u}{\partial t^{2}}=\frac{\partial^{2} u}{\partial x^{2}}+\frac{\partial^{2} u}{\partial y^{2}}+\frac{\partial^{2} u}{\partial z^{2}}-\frac{1}{v^{2}} \frac{\partial^{2} u}{\partial t^{2}}=0
$$

where $u$ is the disturbance and $v$ the propagation velocity. In general holds: $v=f \lambda$. By definition holds: $k \lambda=2 \pi$ and $\omega=2 \pi f$.
In principle, there are two types of waves:

1. Longitudinal waves: for these holds $\vec{k}\|\vec{v}\| \vec{u}$.
2. Transversal waves: for these holds $\vec{k} \| \vec{v} \perp \vec{u}$.

The phase velocity is given by $v_{\mathrm{ph}}=\omega / k$. The group velocity is given by:

$$
v_{\mathrm{g}}=\frac{d \omega}{d k}=v_{\mathrm{ph}}+k \frac{d v_{\mathrm{ph}}}{d k}=v_{\mathrm{ph}}\left(1-\frac{k}{n} \frac{d n}{d k}\right)
$$

where $n$ is the refractive index of the medium. If $v_{\mathrm{ph}}$ does not depend on $\omega$ holds: $v_{\mathrm{ph}}=v_{\mathrm{g}}$. In a dispersive medium it is possible that $v_{\mathrm{g}}>v_{\mathrm{ph}}$ or $v_{\mathrm{g}}<v_{\mathrm{ph}}$, and $v_{\mathrm{g}} \cdot v_{\mathrm{f}}=c^{2}$. If one wants to transfer information with a wave, e.g. by modulation of an EM wave, the information travels with the velocity at with a change in the electromagnetic field propagates. This velocity is often almost equal to the group velocity.
For some media, the propagation velocity follows from:

- Pressure waves in a liquid or gas: $v=\sqrt{\kappa / \varrho}$, where $\kappa$ is the modulus of compression.
- For pressure waves in a gas also holds: $v=\sqrt{\gamma p / \varrho}=\sqrt{\gamma R T / M}$.
- Pressure waves in a thin solid bar with diameter $\ll \lambda: v=\sqrt{E / \varrho}$
- waves in a string: $v=\sqrt{F_{\text {span }} l / m}$
- Surface waves on a liquid: $v=\sqrt{\left(\frac{g \lambda}{2 \pi}+\frac{2 \pi \gamma}{\varrho \lambda}\right) \tanh \left(\frac{2 \pi h}{\lambda}\right)}$
where $h$ is the depth of the liquid and $\gamma$ the surface tension. If $h \ll \lambda$ holds: $v \approx \sqrt{g h}$.


### 5.2 Solutions of the wave equation

### 5.2.1 Plane waves

In $n$ dimensions a harmonic plane wave is defined by:

$$
u(\vec{x}, t)=2^{n} \hat{u} \cos (\omega t) \sum_{i=1}^{n} \sin \left(k_{i} x_{i}\right)
$$

The equation for a harmonic traveling plane wave is: $u(\vec{x}, t)=\hat{u} \cos (\vec{k} \cdot \vec{x} \pm \omega t+\varphi)$
If waves reflect at the end of a spring this will result in a change in phase. A fixed end gives a phase change of $\pi / 2$ to the reflected wave, with boundary condition $u(l)=0$. A lose end gives no change in the phase of the reflected wave, with boundary condition $(\partial u / \partial x)_{l}=0$.

If an observer is moving w.r.t. the wave with a velocity $v_{\text {obs }}$, he will observe a change in frequency: the Doppler effect. This is given by: $\frac{f}{f_{0}}=\frac{v_{\mathrm{f}}-v_{\mathrm{obs}}}{v_{\mathrm{f}}}$.

### 5.2.2 Spherical waves

When the situation is spherical symmetric, the homogeneous wave equation is given by:

$$
\frac{1}{v^{2}} \frac{\partial^{2}(r u)}{\partial t^{2}}-\frac{\partial^{2}(r u)}{\partial r^{2}}=0
$$

with general solution:

$$
u(r, t)=C_{1} \frac{f(r-v t)}{r}+C_{2} \frac{g(r+v t)}{r}
$$

### 5.2.3 Cylindrical waves

When the situation has a cylindrical symmetry, the homogeneous wave equation becomes:

$$
\frac{1}{v^{2}} \frac{\partial^{2} u}{\partial t^{2}}-\frac{1}{r} \frac{\partial}{\partial r}\left(r \frac{\partial u}{\partial r}\right)=0
$$

This is a Bessel equation, with solutions which can be written as Hankel functions. For sufficient large values of $r$ these are approximated by:

$$
u(r, t)=\frac{\hat{u}}{\sqrt{r}} \cos (k(r \pm v t))
$$

### 5.2.4 The general solution in one dimension

Starting point is the equation:

$$
\frac{\partial^{2} u(x, t)}{\partial t^{2}}=\sum_{m=0}^{N}\left(b_{m} \frac{\partial^{m}}{\partial x^{m}}\right) u(x, t)
$$

where $b_{m} \in \mathbb{R}$. Substituting $u(x, t)=A \mathrm{e}^{i(k x-\omega t)}$ gives two solutions $\omega_{j}=\omega_{j}(k)$ as dispersion relations. The general solution is given by:

$$
u(x, t)=\int_{-\infty}^{\infty}\left(a(k) \mathrm{e}^{i\left(k x-\omega_{1}(k) t\right)}+b(k) \mathrm{e}^{i\left(k x-\omega_{2}(k) t\right)}\right) d k
$$

Because in general the frequencies $\omega_{j}$ are non-linear in $k$ there is dispersion and the solution cannot be written any more as a sum of functions depending only on $x \pm v t$ : the wave front transforms.

### 5.3 The stationary phase method

Usually the Fourier integrals of the previous section cannot be calculated exactly. If $\omega_{j}(k) \in \mathbb{R}$ the stationary phase method can be applied. Assuming that $a(k)$ is only a slowly varying function of $k$, one can state that the parts of the $k$-axis where the phase of $k x-\omega(k) t$ changes rapidly will give no net contribution to the integral because the exponent oscillates rapidly there. The only areas contributing significantly to the integral are areas with a stationary phase, determined by $\frac{d}{d k}(k x-\omega(k) t)=0$. Now the following approximation is possible:

$$
\int_{-\infty}^{\infty} a(k) \mathrm{e}^{i(k x-\omega(k) t)} d k \approx \sum_{i=1}^{N} \sqrt{\frac{2 \pi}{\frac{d^{2} \omega\left(k_{i}\right)}{d k_{i}^{2}}}} \exp \left[-i \frac{1}{4} \pi+i\left(k_{i} x-\omega\left(k_{i}\right) t\right)\right]
$$

### 5.4 Green functions for the initial-value problem

This method is preferable if the solutions deviate much from the stationary solutions, like point-like excitations. Starting with the wave equation in one dimension, with $\nabla^{2}=\partial^{2} / \partial x^{2}$ holds: if $Q\left(x, x^{\prime}, t\right)$ is the solution with initial values $Q\left(x, x^{\prime}, 0\right)=\delta\left(x-x^{\prime}\right)$ and $\frac{\partial Q\left(x, x^{\prime}, 0\right)}{\partial t}=0$, and $P\left(x, x^{\prime}, t\right)$ the solution with initial values $P\left(x, x^{\prime}, 0\right)=0$ and $\frac{\partial P\left(x, x^{\prime}, 0\right)}{\partial t}=\delta\left(x-x^{\prime}\right)$, then the solution of the wave equation with arbitrary initial conditions $f(x)=u(x, 0)$ and $g(x)=\frac{\partial u(x, 0)}{\partial t}$ is given by:

$$
u(x, t)=\int_{-\infty}^{\infty} f\left(x^{\prime}\right) Q\left(x, x^{\prime}, t\right) d x^{\prime}+\int_{-\infty}^{\infty} g\left(x^{\prime}\right) P\left(x, x^{\prime}, t\right) d x^{\prime}
$$

$P$ and $Q$ are called the propagators. They are defined by:

$$
\begin{aligned}
Q\left(x, x^{\prime}, t\right) & =\frac{1}{2}\left[\delta\left(x-x^{\prime}-v t\right)+\delta\left(x-x^{\prime}+v t\right)\right] \\
P\left(x, x^{\prime}, t\right) & = \begin{cases}\frac{1}{2 v} & \text { if }\left|x-x^{\prime}\right|<v t \\
0 & \text { if }\left|x-x^{\prime}\right|>v t\end{cases}
\end{aligned}
$$

Further holds the relation: $Q\left(x, x^{\prime}, t\right)=\frac{\partial P\left(x, x^{\prime}, t\right)}{\partial t}$

### 5.5 Waveguides and resonating cavities

The boundary conditions for a perfect conductor can be derived from the Maxwell equations. If $\vec{n}$ is a unit vector $\perp$ the surface, pointed from 1 to 2 , and $\vec{K}$ is a surface current density, than holds:

$$
\begin{array}{ll}
\vec{n} \cdot\left(\vec{D}_{2}-\vec{D}_{1}\right)=\sigma & \vec{n} \times\left(\vec{E}_{2}-\vec{E}_{1}\right)=0 \\
\vec{n} \cdot\left(\vec{B}_{2}-\vec{B}_{1}\right)=0 & \vec{n} \times\left(\vec{H}_{2}-\vec{H}_{1}\right)=\vec{K}
\end{array}
$$

In a waveguide holds because of the cylindrical symmetry: $\vec{E}(\vec{x}, t)=\overrightarrow{\mathcal{E}}(x, y) \mathrm{e}^{i(k z-\omega t)}$ and $\vec{B}(\vec{x}, t)=$ $\overrightarrow{\mathcal{B}}(x, y) \mathrm{e}^{i(k z-\omega t)}$. From this one can now deduce that, if $\mathcal{B}_{z}$ and $\mathcal{E}_{z}$ are not $\equiv 0$ :

$$
\begin{array}{ll}
\mathcal{B}_{x}=\frac{i}{\varepsilon \mu \omega^{2}-k^{2}}\left(k \frac{\partial \mathcal{B}_{z}}{\partial x}-\varepsilon \mu \omega \frac{\partial \mathcal{E}_{z}}{\partial y}\right) & \mathcal{B}_{y}=\frac{i}{\varepsilon \mu \omega^{2}-k^{2}}\left(k \frac{\partial \mathcal{B}_{z}}{\partial y}+\varepsilon \mu \omega \frac{\partial \mathcal{E}_{z}}{\partial x}\right) \\
\mathcal{E}_{x}=\frac{i}{\varepsilon \mu \omega^{2}-k^{2}}\left(k \frac{\partial \mathcal{E}_{z}}{\partial x}+\varepsilon \mu \omega \frac{\partial \mathcal{B}_{z}}{\partial y}\right) & \mathcal{E}_{y}=\frac{i}{\varepsilon \mu \omega^{2}-k^{2}}\left(k \frac{\partial \mathcal{E}_{z}}{\partial y}-\varepsilon \mu \omega \frac{\partial \mathcal{B}_{z}}{\partial x}\right)
\end{array}
$$

Now one can distinguish between three cases:

1. $B_{z} \equiv 0$ : the Transversal Magnetic modes (TM). Boundary condition: $\left.\mathcal{E}_{z}\right|_{\text {surf }}=0$.
2. $E_{z} \equiv 0$ : the Transversal Electric modes (TE). Boundary condition: $\left.\frac{\partial \mathcal{B}_{z}}{\partial n}\right|_{\text {surf }}=0$.

For the TE and TM modes this gives an eigenvalue problem for $\mathcal{E}_{z}$ resp. $\mathcal{B}_{z}$ with boundary conditions:

$$
\left(\frac{\partial^{2}}{\partial x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right) \psi=-\gamma^{2} \psi \text { with eigenvalues } \gamma^{2}:=\varepsilon \mu \omega^{2}-k^{2}
$$

This gives a discrete solution $\psi_{\ell}$ with eigenvalue $\gamma_{\ell}^{2}: k=\sqrt{\varepsilon \mu \omega^{2}-\gamma_{\ell}^{2}}$. For $\omega<\omega_{\ell}, k$ is imaginary and the wave is damped. Therefore, $\omega_{\ell}$ is called the cut-off frequency. In rectangular conductors the following expression can be found for the cut-off frequency for modes $\mathrm{TE}_{m, n}$ of $\mathrm{TM}_{m, n}$ :

$$
\lambda_{\ell}=\frac{2}{\sqrt{(m / a)^{2}+(n / b)^{2}}}
$$

3. $E_{z}$ and $B_{z}$ are zero everywhere: the Transversal electromagnetic mode (TEM). Than holds: $k=$ $\pm \omega \sqrt{\varepsilon \mu}$ and $v_{\mathrm{f}}=v_{\mathrm{g}}$, just as if here were no waveguide. Further $k \in \mathbb{R}$, so there exists no cut-off frequency.

In a rectangular, 3 dimensional resonating cavity with edges $a, b$ and $c$ the possible wave numbers are given by: $k_{x}=\frac{n_{1} \pi}{a}, k_{y}=\frac{n_{2} \pi}{b}, k_{z}=\frac{n_{3} \pi}{c}$ This results in the possible frequencies $f=v k / 2 \pi$ in the cavity:

$$
f=\frac{v}{2} \sqrt{\frac{n_{x}^{2}}{a^{2}}+\frac{n_{y}^{2}}{b^{2}}+\frac{n_{z}^{2}}{c^{2}}}
$$

For a cubic cavity, with $a=b=c$, the possible number of oscillating modes $N_{\mathrm{L}}$ for longitudinal waves is given by:

$$
N_{\mathrm{L}}=\frac{4 \pi a^{3} f^{3}}{3 v^{3}}
$$

Because transversal waves have two possible polarizations holds for them: $N_{\mathrm{T}}=2 N_{\mathrm{L}}$.

### 5.6 Non-linear wave equations

The Van der Pol equation is given by:

$$
\frac{d^{2} x}{d t^{2}}-\varepsilon \omega_{0}\left(1-\beta x^{2}\right) \frac{d x}{d t}+\omega_{0}^{2} x=0
$$

$\beta x^{2}$ can be ignored for very small values of the amplitude. Substitution of $x \sim \mathrm{e}^{i \omega t}$ gives: $\omega=\frac{1}{2} \omega_{0}(i \varepsilon \pm$ $2 \sqrt{1-\frac{1}{2} \varepsilon^{2}}$ ). The lowest-order instabilities grow as $\frac{1}{2} \varepsilon \omega_{0}$. While $x$ is growing, the 2 nd term becomes larger and diminishes the growth. Oscillations on a time scale $\sim \omega_{0}^{-1}$ can exist. If $x$ is expanded as $x=x^{(0)}+$ $\varepsilon x^{(1)}+\varepsilon^{2} x^{(2)}+\cdots$ and this is substituted one obtains, besides periodic, secular terms $\sim \varepsilon t$. If it is assumed that there exist timescales $\tau_{n}, 0 \leq \tau \leq N$ with $\partial \tau_{n} / \partial t=\varepsilon^{n}$ and if the secular terms are put 0 one obtains:

$$
\frac{d}{d t}\left\{\frac{1}{2}\left(\frac{d x}{d t}\right)^{2}+\frac{1}{2} \omega_{0}^{2} x^{2}\right\}=\varepsilon \omega_{0}\left(1-\beta x^{2}\right)\left(\frac{d x}{d t}\right)^{2}
$$

This is an energy equation. Energy is conserved if the left-hand side is 0 . If $x^{2}>1 / \beta$, the right-hand side changes sign and an increase in energy changes into a decrease of energy. This mechanism limits the growth of oscillations.

The Korteweg-De Vries equation is given by:

$$
\frac{\partial u}{\partial t}+\frac{\partial u}{\partial x}-\underbrace{a u \frac{\partial u}{\partial x}}_{\text {non-lin }}+\underbrace{b^{2} \frac{\partial^{3} u}{\partial x^{3}}}_{\text {dispersive }}=0
$$

This equation is for example a model for ion-acoustic waves in a plasma. For this equation, soliton solutions of the following form exist:

$$
u(x-c t)=\frac{-d}{\cosh ^{2}(e(x-c t))}
$$

with $c=1+\frac{1}{3} a d$ and $e^{2}=a d /\left(12 b^{2}\right)$.

## Chapter 6

## Optics

### 6.1 The bending of light

For the refraction at a surface holds: $n_{i} \sin \left(\theta_{i}\right)=n_{t} \sin \left(\theta_{t}\right)$ where $n$ is the refractive index of the material. Snell's law is:

$$
\frac{n_{2}}{n_{1}}=\frac{\lambda_{1}}{\lambda_{2}}=\frac{v_{1}}{v_{2}}
$$

If $\Delta n \leq 1$, the change in phase of the light is $\Delta \varphi=0$, if $\Delta n>1$ holds: $\Delta \varphi=\pi$. The refraction of light in a material is caused by scattering from atoms. This is described by:

$$
n^{2}=1+\frac{n_{\mathrm{e}} e^{2}}{\varepsilon_{0} m} \sum_{j} \frac{f_{j}}{\omega_{0, j}^{2}-\omega^{2}-i \delta \omega}
$$

where $n_{\mathrm{e}}$ is the electron density and $f_{j}$ the oscillator strength, for which holds: $\sum_{j} f_{j}=1$. From this follows that $v_{\mathrm{g}}=c /\left(1+\left(n_{\mathrm{e}} e^{2} / 2 \varepsilon_{0} m \omega^{2}\right)\right)$. From this the equation of Cauchy can be derived: $n=a_{0}+a_{1} / \lambda^{2}$. More general, it is possible to expand $n$ as: $n=\sum_{k=0}^{n} \frac{a_{k}}{\lambda^{2 k}}$.

For an electromagnetic wave in general holds: $n=\sqrt{\varepsilon_{\mathrm{r}} \mu_{\mathrm{r}}}$.
The path, followed by a light ray in material can be found from Fermat's principle:

$$
\delta \int_{1}^{2} d t=\delta \int_{1}^{2} \frac{n(s)}{c} d s=0 \Rightarrow \delta \int_{1}^{2} n(s) d s=0
$$

### 6.2 Paraxial geometrical optics

### 6.2.1 Lenses

The Gaussian lens formula can be deduced from Fermat's principle with the approximations $\cos \varphi=1$ and $\sin \varphi=\varphi$. For the refraction at a spherical surface with radius $R$ holds:

$$
\frac{n_{1}}{v}-\frac{n_{2}}{b}=\frac{n_{1}-n_{2}}{R}
$$

where $|v|$ is the distance of the object and $|b|$ the distance of the image. Applying this twice results in:

$$
\frac{1}{f}=\left(n_{1}-1\right)\left(\frac{1}{R_{2}}-\frac{1}{R_{1}}\right)
$$

where $n_{1}$ is the refractive index of the lens, $f$ is the focal length and $R_{1}$ and $R_{2}$ are the curvature radii of both surfaces. For a double concave lens holds $R_{1}<0, R_{2}>0$, for a double convex lens holds $R_{1}>0$ and $R_{2}<0$. Further holds:

$$
\frac{1}{f}=\frac{1}{v}-\frac{1}{b}
$$

$D:=1 / f$ is called the dioptric power of a lens. For a lens with thickness $d$ and diameter $D$ holds to a good approximation: $1 / f=8(n-1) d / D^{2}$. For two lenses placed on a line with distance $d$ holds:

$$
\frac{1}{f}=\frac{1}{f_{1}}+\frac{1}{f_{2}}-\frac{d}{f_{1} f_{2}}
$$

In these equations the following signs are being used for refraction at a spherical surface, as is seen by an incoming light ray:

| Quantity | + | - |
| :---: | :--- | :--- |
| $R$ | Concave surface | Convex surface |
| $f$ | Converging lens | Diverging lens |
| $v$ | Real object | Virtual object |
| $b$ | Virtual image | Real image |

### 6.2.2 Mirrors

For images of mirrors holds:

$$
\frac{1}{f}=\frac{1}{v}+\frac{1}{b}=\frac{2}{R}+\frac{h^{2}}{2}\left(\frac{1}{R}-\frac{1}{v}\right)^{2}
$$

where $h$ is the perpendicular distance from the point the light ray hits the mirror to the optical axis. Spherical aberration can be reduced by not using spherical mirrors. A parabolical mirror has no spherical aberration for light rays parallel with the optical axis and is therefore often used for telescopes. The used signs are:

| Quantity | + | - |
| :---: | :---: | :---: |
| $R$ | Concave mirror | Convex mirror |
| $f$ | Concave mirror | Convex mirror |
| $v$ | Real object | Virtual object |
| $b$ | Real image | Virtual image |

### 6.2.3 Principal planes

The nodal points N of a lens are defined by the figure on the right. If the lens is surrounded by the same medium on both sides, the nodal points are the same as the principal points H . The plane $\perp$ the optical axis through the principal points is called the principal plane. If the lens is described by a matrix $m_{i j}$ than for the distances $h_{1}$ and $h_{2}$ to the boundary of the lens holds:

$$
h_{1}=n \frac{m_{11}-1}{m_{12}}, \quad h_{2}=n \frac{m_{22}-1}{m_{12}}
$$

### 6.2.4 Magnification

The linear magnification is defined by: $N=-\frac{b}{v}$
The angular magnification is defined by: $N_{\alpha}=-\frac{\alpha_{\text {syst }}}{\alpha_{\text {none }}}$
where $\alpha_{\text {sys }}$ is the size of the retinal image with the optical system and $\alpha_{\text {none }}$ the size of the retinal image without the system. Further holds: $N \cdot N_{\alpha}=1$. For a telescope holds: $N=f_{\text {objective }} / f_{\text {ocular }}$. The $f$-number is defined by $f / D_{\text {objective }}$.

### 6.3 Matrix methods

A light ray can be described by a vector $(n \alpha, y)$ with $\alpha$ the angle with the optical axis and $y$ the distance to the optical axis. The change of a light ray interacting with an optical system can be obtained using a matrix multiplication:

$$
\binom{n_{2} \alpha_{2}}{y_{2}}=M\binom{n_{1} \alpha_{1}}{y_{1}}
$$

where $\operatorname{Tr}(M)=1 . M$ is a product of elementary matrices. These are:

1. Transfer along length $l: M_{\mathrm{R}}=\left(\begin{array}{cc}1 & 0 \\ l / n & 1\end{array}\right)$
2. Refraction at a surface with dioptric power $D: M_{\mathrm{T}}=\left(\begin{array}{cc}1 & -D \\ 0 & 1\end{array}\right)$

### 6.4 Aberrations

Lenses usually do not give a perfect image. Some causes are:

1. Chromatic aberration is caused by the fact that $n=n(\lambda)$. This can be partially corrected with a lens which is composed of more lenses with different functions $n_{i}(\lambda)$. Using $N$ lenses makes it possible to obtain the same $f$ for $N$ wavelengths.
2. Spherical aberration is caused by second-order effects which are usually ignored; a spherical surface does not make a perfect lens. Incomming rays far from the optical axis will more bent.
3. Coma is caused by the fact that the principal planes of a lens are only flat near the principal axis. Further away of the optical axis they are curved. This curvature can be both positive or negative.
4. Astigmatism: from each point of an object not on the optical axis the image is an ellipse because the thickness of the lens is not the same everywhere.
5. Field curvature can be corrected by the human eye.
6. Distorsion gives abberations near the edges of the image. This can be corrected with a combination of positive and negative lenses.

### 6.5 Reflection and transmission

If an electromagnetic wave hits a transparent medium part of the wave will reflect at the same angle as the incident angle, and a part will be refracted at an angle according to Snell's law. It makes a difference whether the $\vec{E}$ field of the wave is $\perp$ or $\|$ w.r.t. the surface. When the coefficients of reflection $r$ and transmission $t$ are defined as:

$$
r_{\|} \equiv\left(\frac{E_{0 r}}{E_{0 i}}\right)_{\|}, r_{\perp} \equiv\left(\frac{E_{0 r}}{E_{0 i}}\right)_{\perp}, t_{\|} \equiv\left(\frac{E_{0 t}}{E_{0 i}}\right)_{\|}, t_{\perp} \equiv\left(\frac{E_{0 t}}{E_{0 i}}\right)_{\perp}
$$

where $E_{0 r}$ is the reflected amplitude and $E_{0 t}$ the transmitted amplitude. Then the Fresnel equations are:

$$
\begin{gathered}
r_{\|}=\frac{\tan \left(\theta_{i}-\theta_{t}\right)}{\tan \left(\theta_{i}+\theta_{t}\right)}, \quad r_{\perp}=\frac{\sin \left(\theta_{t}-\theta_{i}\right)}{\sin \left(\theta_{t}+\theta_{i}\right)} \\
t_{\|}=\frac{2 \sin \left(\theta_{t}\right) \cos \left(\theta_{i}\right)}{\sin \left(\theta_{t}+\theta_{i}\right) \cos \left(\theta_{t}-\theta_{i}\right)}, \quad t_{\perp}=\frac{2 \sin \left(\theta_{t}\right) \cos \left(\theta_{i}\right)}{\sin \left(\theta_{t}+\theta_{i}\right)}
\end{gathered}
$$

The following holds: $t_{\perp}-r_{\perp}=1$ and $t_{\|}+r_{\|}=1$. If the coefficient of reflection $R$ and transmission $T$ are defined as (with $\theta_{i}=\theta_{r}$ ):

$$
R \equiv \frac{I_{r}}{I_{i}} \quad \text { and } \quad T \equiv \frac{I_{t} \cos \left(\theta_{t}\right)}{I_{i} \cos \left(\theta_{i}\right)}
$$

with $I=\langle | \vec{S}| \rangle$ it follows: $R+T=1$. A special case is $r_{\|}=0$. This happens if the angle between the reflected and transmitted rays is $90^{\circ}$. From Snell's law it then follows: $\tan \left(\theta_{i}\right)=n$. This angle is called Brewster's angle. The situation with $r_{\perp}=0$ is not possible.

### 6.6 Polarization

The polarization is defined as: $P=\frac{I_{\mathrm{p}}}{I_{\mathrm{p}}+I_{\mathrm{u}}}=\frac{I_{\text {max }}-I_{\text {min }}}{I_{\max }+I_{\text {min }}}$
where the intensity of the polarized light is given by $I_{\mathrm{p}}$ and the intensity of the unpolarized light is given by $I_{\mathrm{u}} . I_{\max }$ and $I_{\min }$ are the maximum and minimum intensities when the light passes a polarizer. If polarized light passes through a polarizer Malus law applies: $I(\theta)=I(0) \cos ^{2}(\theta)$ where $\theta$ is the angle of the polarizer.
The state of a light ray can be described by the Stokes-parameters: start with 4 filters which each transmits half the intensity. The first is independent of the polarization, the second and third are linear polarizers with the transmission axes horizontal and at $+45^{\circ}$, while the fourth is a circular polarizer which is opaque for $L$-states. Then holds $S_{1}=2 I_{1}, S_{2}=2 I_{2}-2 I_{1}, S_{3}=2 I_{3}-2 I_{1}$ and $S_{4}=2 I_{4}-2 I_{1}$.
The state of a polarized light ray can also be described by the Jones vector:

$$
\vec{E}=\binom{E_{0 x} \mathrm{e}^{i \varphi_{x}}}{E_{0 y} \mathrm{e}^{i \varphi_{y}}}
$$

For the horizontal $P$-state holds: $\vec{E}=(1,0)$, for the vertical $P$-state $\vec{E}=(0,1)$, the $R$-state is given by $\vec{E}=\frac{1}{2} \sqrt{2}(1,-i)$ and the $L$-state by $\vec{E}=\frac{1}{2} \sqrt{2}(1, i)$. The change in state of a light beam after passage of optical equipment can be described as $\vec{E}_{2}=M \cdot \vec{E}_{1}$. For some types of optical equipment the Jones matrix $M$ is given by:

| Horizontal linear polarizer: | $\left(\begin{array}{cc}1 & 0 \\ 0 & 0\end{array}\right)$ |
| :--- | ---: |
| Vertical linear polarizer: | $\left(\begin{array}{ll}0 & 0 \\ 0 & 1\end{array}\right)$ |
| Linear polarizer at $+45^{\circ}$ | $\frac{1}{2}\left(\begin{array}{cc}1 & 1 \\ 1 & 1\end{array}\right)$ |
| Lineair polarizer at $-45^{\circ}$ | $\frac{1}{2}\left(\begin{array}{cc}1 & -1 \\ -1 & 1\end{array}\right)$ |
| $\frac{1}{4}-\lambda$ plate, fast axis vertical | $\mathrm{e}^{i \pi / 4}\left(\begin{array}{cc}1 & 0 \\ 0 & -i\end{array}\right)$ |
| $\frac{1}{4}-\lambda$ plate, fast axis horizontal | $\mathrm{e}^{i \pi / 4}\left(\begin{array}{cc}1 & 0 \\ 0 & i\end{array}\right)$ |
| Homogene circular polarizor right | $\frac{1}{2}\left(\begin{array}{cc}1 & i \\ -i & 1\end{array}\right)$ |
| Homogene circular polarizer left | $\frac{1}{2}\left(\begin{array}{cc}1 & -i \\ i & 1\end{array}\right)$ |

### 6.7 Prisms and dispersion

A light ray passing through a prism is refracted twice and aquires a deviation from its original direction $\delta=\theta_{i}+\theta_{i^{\prime}}+\alpha$ w.r.t. the incident direction, where $\alpha$ is the apex angle, $\theta_{i}$ is the angle between the incident angle and a line perpendicular to the surface and $\theta_{i^{\prime}}$ is the angle between the ray leaving the prism and a line perpendicular to the surface. When $\theta_{i}$ varies there is an angle for which $\delta$ becomes minimal. For the refractive index of the prism now holds:

$$
n=\frac{\sin \left(\frac{1}{2}\left(\delta_{\min }+\alpha\right)\right)}{\sin \left(\frac{1}{2} \alpha\right)}
$$

The dispersion of a prism is defined by:

$$
D=\frac{d \delta}{d \lambda}=\frac{d \delta}{d n} \frac{d n}{d \lambda}
$$

where the first factor depends on the shape and the second on the composition of the prism. For the first factor follows:

$$
\frac{d \delta}{d n}=\frac{2 \sin \left(\frac{1}{2} \alpha\right)}{\cos \left(\frac{1}{2}\left(\delta_{\min }+\alpha\right)\right)}
$$

For visible light usually holds $d n / d \lambda<0$ : shorter wavelengths are stronger bent than longer. The refractive index in this area can usually be approximated by Cauchy's formula.

### 6.8 Diffraction

Fraunhofer diffraction occurs far away from the source(s). The Fraunhofer diffraction of light passing through multiple slits is described by:

$$
\frac{I(\theta)}{I_{0}}=\left(\frac{\sin (u)}{u}\right)^{2} \cdot\left(\frac{\sin (N v)}{\sin (v)}\right)^{2}
$$

where $u=\pi b \sin (\theta) / \lambda, v=\pi d \sin (\theta) / \lambda . N$ is the number of slits, $b$ the width of a slit and $d$ the distance between the slits. The maxima in intensity are given by $d \sin (\theta)=k \lambda$.
The diffraction through a spherical aperture with radius $a$ is described by:

$$
\frac{I(\theta)}{I_{0}}=\left(\frac{J_{1}(k a \sin (\theta))}{k a \sin (\theta)}\right)^{2}
$$

The diffraction pattern of a rectangular aperture at distance $R$ with length $a$ in the $x$-direction and $b$ in the $y$-direction is described by:

$$
\frac{I(x, y)}{I_{0}}=\left(\frac{\sin \left(\alpha^{\prime}\right)}{\alpha^{\prime}}\right)^{2}\left(\frac{\sin \left(\beta^{\prime}\right)}{\beta^{\prime}}\right)^{2}
$$

where $\alpha^{\prime}=k a x / 2 R$ and $\beta^{\prime}=k b y / 2 R$.
When X rays are diffracted at a crystal holds for the position of the maxima in intensity Bragg's relation: $2 d \sin (\theta)=n \lambda$ where $d$ is the distance between the crystal layers.
Close at the source the Fraunhofermodel is invalid because it ignores the angle-dependence of the reflected waves. This is described by the obliquity or inclination factor, which describes the directionality of the secondary emissions: $E(\theta)=\frac{1}{2} E_{0}(1+\cos (\theta))$ where $\theta$ is the angle w.r.t. the optical axis.
Diffraction limits the resolution of a system. This is the minimum angle $\Delta \theta_{\text {min }}$ between two incident rays coming from points far away for which their refraction patterns can be detected separately. For a circular slit holds: $\Delta \theta_{\min }=1.22 \lambda / D$ where $D$ is the diameter of the slit.
For a grating holds: $\Delta \theta_{\min }=2 \lambda /\left(N a \cos \left(\theta_{m}\right)\right)$ where $a$ is the distance between two peaks and $N$ the number of peaks. The minimum difference between two wavelengths that gives a separated diffraction pattern in a multiple slit geometry is given by $\Delta \lambda / \lambda=n N$ where $N$ is the number of lines and $n$ the order of the pattern.

### 6.9 Special optical effects

- Birefringe and dichroism. $\vec{D}$ is not parallel with $\vec{E}$ if the polarizability $\vec{P}$ of a material is not equal in all directions. There are at least 3 directions, the principal axes, in which they are parallel. This results in 3 refractive indices $n_{i}$ which can be used to construct Fresnel's ellipsoid. In case $n_{2}=n_{3} \neq n_{1}$, which happens e.g. at trigonal, hexagonal and tetragonal crystals there is one optical axis in the direction of $n_{1}$. Incident light rays can now be split up in two parts: the ordinary wave is linear polarized $\perp$ the plane through the transmission direction and the optical axis. The extraordinary wave is linear polarized
in the plane through the transmission direction and the optical axis. Dichroism is caused by a different absorption of the ordinary and extraordinary wave in some materials. Double images occur when the incident ray makes an angle with the optical axis: the extraordinary wave will refract, the ordinary will not.
- Retarders: waveplates and compensators. Incident light will have a phase shift of $\Delta \varphi=2 \pi d\left(\mid n_{0}-\right.$ $\left.n_{\mathrm{e}} \mid\right) / \lambda_{0}$ if an uniaxial crystal is cut in such a way that the optical axis is parallel with the front and back plane. Here, $\lambda_{0}$ is the wavelength in vacuum and $n_{0}$ and $n_{\mathrm{e}}$ the refractive indices for the ordinary and extraordinary wave. For a quarter-wave plate holds: $\Delta \varphi=\pi / 2$.
- The Kerr-effect: isotropic, transparent materials can become birefringent when placed in an electric field. In that case, the optical axis is parallel to $\vec{E}$. The difference in refractive index in the two directions is given by: $\Delta n=\lambda_{0} K E^{2}$, where $K$ is the Kerr constant of the material. If the electrodes have an effective length $\ell$ and are separated by a distance $d$, the retardation is given by: $\Delta \varphi=2 \pi K \ell V^{2} / d^{2}$, where $V$ is the applied voltage.
- The Pockels or linear electro-optical effect can occur in 20 (from a total of 32) crystal symmetry classes, namely those without a centre of symmetry. These crystals are also piezoelectric: their polarization changes when a pressure is applied and vice versa: $\vec{P}=p d+\varepsilon_{0} \chi \vec{E}$. The retardation in a Pockels cell is $\Delta \varphi=2 \pi n_{0}^{3} r_{63} V / \lambda_{0}$ where $r_{63}$ is the 6-3 element of the electro-optic tensor.
- The Faraday effect: the polarization of light passing through material with length $d$ and to which a magnetic field is applied in the propagation direction is rotated by an angle $\beta=\mathcal{V} B d$ where $\mathcal{V}$ is the Verdet constant.
- C̆erenkov radiation arises when a charged particle with $v_{q}>v_{\mathrm{f}}$ arrives. The radiation is emitted within a cone with an apex angle $\alpha$ with $\sin (\alpha)=c / c_{\text {medium }}=c / n v_{q}$.


### 6.10 The Fabry-Perot interferometer

For a Fabry-Perot interferometer holds in general: $T+R+A=1$ where $T$ is the transmission factor, $R$ the reflection factor and $A$ the absorption factor. If $F$ is given by $F=4 R /(1-R)^{2}$ it follows for the intensity distribution:

$$
\frac{I_{t}}{I_{i}}=\left[1-\frac{A}{1-R}\right]^{2} \frac{1}{1+F \sin ^{2}(\theta)}
$$

The term $\left[1+F \sin ^{2}(\theta)\right]^{-1}:=\mathcal{A}(\theta)$ is called the Airy function.


The width of the peaks at half height is given by $\gamma=4 / \sqrt{F}$. The finesse $\mathcal{F}$ is defined as $\mathcal{F}=\frac{1}{2} \pi \sqrt{F}$. The maximum resolution is then given by $\Delta f_{\min }=c / 2 n d \mathcal{F}$.

## Chapter 7

## Statistical physics

### 7.1 Degrees of freedom

A molecule consisting of $n$ atoms has $s=3 n$ degrees of freedom. There are 3 translational degrees of freedom, a linear molecule has $s=3 n-5$ vibrational degrees of freedom and a non-linear molecule $s=3 n-6$. A linear molecule has 2 rotational degrees of freedom and a non-linear molecule 3 .

Because vibrational degrees of freedom account for both kinetic and potential energy they count double. So, for linear molecules this results in a total of $s=6 n-5$. For non-linear molecules this gives $s=6 n-6$. The average energy of a molecule in thermodynamic equilibrium is $\left\langle E_{\text {tot }}\right\rangle=\frac{1}{2} s k T$. Each degree of freedom of a molecule has in principle the same energy: the principle of equipartition.

The rotational and vibrational energy of a molecule are:

$$
W_{\mathrm{rot}}=\frac{\hbar^{2}}{2 I} l(l+1)=B l(l+1), \quad W_{\mathrm{vib}}=\left(v+\frac{1}{2}\right) \hbar \omega_{0}
$$

The vibrational levels are excited if $k T \approx \hbar \omega$, the rotational levels of a hetronuclear molecule are excited if $k T \approx 2 B$. For homonuclear molecules additional selection rules apply so the rotational levels are well coupled if $k T \approx 6 B$.

### 7.2 The energy distribution function

The general form of the equilibrium velocity distribution function is
$P\left(v_{x}, v_{y}, v_{z}\right) d v_{x} d v_{y} d v_{z}=P\left(v_{x}\right) d v_{x} \cdot P\left(v_{y}\right) d v_{y} \cdot P\left(v_{z}\right) d v_{z}$ with

$$
P\left(v_{i}\right) d v_{i}=\frac{1}{\alpha \sqrt{\pi}} \exp \left(-\frac{v_{i}^{2}}{\alpha^{2}}\right) d v_{i}
$$

where $\alpha=\sqrt{2 k T / m}$ is the most probable velocity of a particle. The average velocity is given by $\langle v\rangle=$ $2 \alpha / \sqrt{\pi}$, and $\left\langle v^{2}\right\rangle=\frac{3}{2} \alpha^{2}$. The distribution as a function of the absolute value of the velocity is given by:

$$
\frac{d N}{d v}=\frac{4 N}{\alpha^{3} \sqrt{\pi}} v^{2} \exp \left(-\frac{m v^{2}}{2 k T}\right)
$$

The general form of the energy distribution function then becomes:

$$
P(E) d E=\frac{c(s)}{k T}\left(\frac{E}{k T}\right)^{\frac{1}{2} s-1} \exp \left(-\frac{E}{k T}\right) d E
$$

where $c(s)$ is a normalization constant, given by:

1. Even $s: s=2 l: c(s)=\frac{1}{(l-1)!}$
2. Odd $s: s=2 l+1: c(s)=\frac{2^{l}}{\sqrt{\pi}(2 l-1)!!}$

### 7.3 Pressure on a wall

The number of molecules that collides with a wall with surface $A$ within a time $\tau$ is given by:

$$
\iiint d^{3} N=\int_{0}^{\infty} \int_{0}^{\pi} \int_{0}^{2 \pi} n A v \tau \cos (\theta) P(v, \theta, \varphi) d v d \theta d \varphi
$$

From this follows for the particle flux on the wall: $\Phi=\frac{1}{4} n\langle v\rangle$. For the pressure on the wall then follows:

$$
d^{3} p=\frac{2 m v \cos (\theta) d^{3} N}{A \tau}, \text { so } p=\frac{2}{3} n\langle E\rangle
$$

### 7.4 The equation of state

If intermolecular forces and the volume of the molecules can be neglected then for gases from $p=\frac{2}{3} n\langle E\rangle$ and $\langle E\rangle=\frac{3}{2} k T$ can be derived:

$$
p V=n_{s} R T=\frac{1}{3} N m\left\langle v^{2}\right\rangle
$$

Here, $n_{s}$ is the number of moles particles and $N$ is the total number of particles within volume $V$. If the own volume and the intermolecular forces cannot be neglected the Van der Waals equation can be derived:

$$
\left(p+\frac{a n_{s}^{2}}{V^{2}}\right)\left(V-b n_{s}\right)=n_{s} R T
$$

There is an isotherme with a horizontal point of inflection. In the Van der Waals equation this corresponds with the critical temperature, pressure and volume of the gas. This is the upper limit of the area of coexistence between liquid and vapor. From $d p / d V=0$ and $d^{2} p / d V^{2}=0$ follows:

$$
T_{\text {cr }}=\frac{8 a}{27 b R}, \quad p_{\text {cr }}=\frac{a}{27 b^{2}}, \quad V_{\text {cr }}=3 b n_{s}
$$

For the critical point holds: $p_{\text {cr }} V_{m, \mathrm{cr}} / R T_{\mathrm{cr}}=\frac{3}{8}$, which differs from the value of 1 which follows from the general gas law.

Scaled on the critical quantities, with $p^{*}:=p / p_{\text {cr }}, T^{*}=T / T_{\text {cr }}$ and $V_{m}^{*}=V_{m} / V_{m, \text { cr }}$ with $V_{m}:=V / n_{s}$ holds:

$$
\left(p^{*}+\frac{3}{\left(V_{m}^{*}\right)^{2}}\right)\left(V_{m}^{*}-\frac{1}{3}\right)=\frac{8}{3} T^{*}
$$

Gases behave the same for equal values of the reduced quantities: the law of the corresponding states. A virial expansion is used for even more accurate views:

$$
p\left(T, V_{m}\right)=R T\left(\frac{1}{V_{m}}+\frac{B(T)}{V_{m}^{2}}+\frac{C(T)}{V_{m}^{3}}+\cdots\right)
$$

The Boyle temperature $T_{\mathrm{B}}$ is the temperature for which the 2 nd virial coefficient is 0 . In a Van der Waals gas, this happens at $T_{\mathrm{B}}=a / R b$. The inversion temperature $T_{\mathrm{i}}=2 T_{\mathrm{B}}$.

The equation of state for solids and liquids is given by:

$$
\frac{V}{V_{0}}=1+\gamma_{p} \Delta T-\kappa_{T} \Delta p=1+\frac{1}{V}\left(\frac{\partial V}{\partial T}\right)_{p} \Delta T+\frac{1}{V}\left(\frac{\partial V}{\partial p}\right)_{T} \Delta p
$$

### 7.5 Collisions between molecules

The collision probability of a particle in a gas that is translated over a distance $d x$ is given by $n \sigma d x$, where $\sigma$ is the cross section. The mean free path is given by $\ell=\frac{v_{1}}{n u \sigma}$ with $u=\sqrt{v_{1}^{2}+v_{2}^{2}}$ the relative velocity between the particles. If $m_{1} \ll m_{2}$ holds: $\frac{u}{v_{1}}=\sqrt{1+\frac{m_{1}}{m_{2}}}$, so $\ell=\frac{1}{n \sigma}$. If $m_{1}=m_{2}$ holds: $\ell=\frac{1}{n \sigma \sqrt{2}}$. This means that the average time between two collisions is given by $\tau=\frac{1}{n \sigma v}$. If the molecules are approximated by hard spheres the cross section is: $\sigma=\frac{1}{4} \pi\left(D_{1}^{2}+D_{2}^{2}\right)$. The average distance between two molecules is $0.55 n^{-1 / 3}$. Collisions between molecules and small particles in a solution result in the Brownian motion. For the average motion of a particle with radius $R$ can be derived: $\left\langle x_{i}^{2}\right\rangle=\frac{1}{3}\left\langle r^{2}\right\rangle=k T t / 3 \pi \eta R$.
A gas is called a Knudsen gas if $\ell \gg$ the dimensions of the gas, something that can easily occur at low pressures. The equilibrium condition for a vessel which has a hole with surface $A$ in it for which holds that $\ell \gg \sqrt{A / \pi}$ is: $n_{1} \sqrt{T_{1}}=n_{2} \sqrt{T_{2}}$. Together with the general gas law follows: $p_{1} / \sqrt{T_{1}}=p_{2} / \sqrt{T_{2}}$.

If two plates move along each other at a distance $d$ with velocity $w_{x}$ the viscosity $\eta$ is given by: $F_{x}=\eta \frac{A w_{x}}{d}$. The velocity profile between the plates is in that case given by $w(z)=z w_{x} / d$. It can be derived that $\eta=$ $\frac{1}{3} \varrho \ell\langle v\rangle$ where $v$ is the thermal velocity.
The heat conductance in a non-moving gas is described by: $\frac{d Q}{d t}=\kappa A\left(\frac{T_{2}-T_{1}}{d}\right)$, which results in a temperature profile $T(z)=T_{1}+z\left(T_{2}-T_{1}\right) / d$. It can be derived that $\kappa=\frac{1}{3} C_{m V} n \ell\langle v\rangle / N_{\mathrm{A}}$. Also holds: $\kappa=C_{V} \eta$. A better expression for $\kappa$ can be obtained with the Eucken correction: $\kappa=\left(1+9 R / 4 c_{m V}\right) C_{V} \cdot \eta$ with an error $<5 \%$.

### 7.6 Interaction between molecules

For dipole interaction between molecules can be derived that $U \sim-1 / r^{6}$. If the distance between two molecules approaches the molecular diameter $D$ a repulsing force between the electron clouds appears. This force can be described by $U_{\text {rep }} \sim \exp (-\gamma r)$ or $V_{\text {rep }}=+C_{s} / r^{s}$ with $12 \leq s \leq 20$. This results in the Lennard-Jones potential for intermolecular forces:

$$
U_{\mathrm{LJ}}=4 \epsilon\left[\left(\frac{D}{r}\right)^{12}-\left(\frac{D}{r}\right)^{6}\right]
$$

with a minimum $\epsilon$ at $r=r_{\mathrm{m}}$. The following holds: $D \approx 0.89 r_{\mathrm{m}}$. For the Van der Waals coefficients $a$ and $b$ and the critical quantities holds: $a=5.275 N_{\mathrm{A}}^{2} D^{3} \epsilon, b=1.3 N_{\mathrm{A}} D^{3}, k T_{\mathrm{kr}}=1.2 \epsilon$ and $V_{\mathrm{m}, \mathrm{kr}}=3.9 N_{\mathrm{A}} D^{3}$.
A more simple model for intermolecular forces assumes a potential $U(r)=\infty$ for $r<D, U(r)=U_{\mathrm{LJ}}$ for $D \leq r \leq 3 D$ and $U(r)=0$ for $r \geq 3 D$. This gives for the potential energy of one molecule: $E_{\mathrm{pot}}=$ $\int_{D}^{3 \bar{D}} U(r) F(r) d r$.
with $F(r)$ the spatial distribution function in spherical coordinates, which for a homogeneous distribution is given by: $F(r) d r=4 n \pi r^{2} d r$.

Some useful mathematical relations are:

$$
\int_{0}^{\infty} x^{n} \mathrm{e}^{-x} d x=n!, \quad \int_{0}^{\infty} x^{2 n} \mathrm{e}^{-x^{2}} d x=\frac{(2 n)!\sqrt{\pi}}{n!2^{2 n+1}}, \quad \int_{0}^{\infty} x^{2 n+1} \mathrm{e}^{-x^{2}} d x=\frac{1}{2} n!
$$

## Chapter 8

## Thermodynamics

### 8.1 Mathematical introduction

If there exists a relation $f(x, y, z)=0$ between 3 variables, one can write: $x=x(y, z), y=y(x, z)$ and $z=z(x, y)$. The total differential $d z$ of $z$ is than given by:

$$
d z=\left(\frac{\partial z}{\partial x}\right)_{y} d x+\left(\frac{\partial z}{\partial y}\right)_{x} d y
$$

By writing this also for $d x$ and $d y$ it can be obtained that

$$
\left(\frac{\partial x}{\partial y}\right)_{z} \cdot\left(\frac{\partial y}{\partial z}\right)_{x} \cdot\left(\frac{\partial z}{\partial x}\right)_{y}=-1
$$

Because $d z$ is a total differential holds $\oint d z=0$.
A homogeneous function of degree $m$ obeys: $\varepsilon^{m} F(x, y, z)=F(\varepsilon x, \varepsilon y, \varepsilon z)$. For such a function Euler's theorem applies:

$$
m F(x, y, z)=x \frac{\partial F}{\partial x}+y \frac{\partial F}{\partial y}+z \frac{\partial F}{\partial z}
$$

### 8.2 Definitions

- The isochoric pressure coefficient: $\beta_{V}=\frac{1}{p}\left(\frac{\partial p}{\partial T}\right)_{V}$
- The isothermal compressibility: $\kappa_{T}=-\frac{1}{V}\left(\frac{\partial V}{\partial p}\right)_{T}$
- The isobaric volume coefficient: $\gamma_{p}=\frac{1}{V}\left(\frac{\partial V}{\partial T}\right)_{p}$
- The adiabatic compressibility: $\kappa_{S}=-\frac{1}{V}\left(\frac{\partial V}{\partial p}\right)_{S}$

For an ideal gas follows: $\gamma_{p}=1 / T, \kappa_{T}=1 / p$ and $\beta_{V}=-1 / V$.

### 8.3 Thermal heat capacity

- The specific heat at constant $X$ is: $C_{X}=T\left(\frac{\partial S}{\partial T}\right)_{X}$
- The specific heat at constant pressure: $C_{p}=\left(\frac{\partial H}{\partial T}\right)_{p}$
- The specific heat at constant volume: $C_{V}=\left(\frac{\partial U}{\partial T}\right)_{V}$

For an ideal gas holds: $C_{m p}-C_{m V}=R$. Further, if the temperature is high enough to thermalize all internal rotational and vibrational degrees of freedom, holds: $C_{V}=\frac{1}{2} s R$. Hence $C_{p}=\frac{1}{2}(s+2) R$. For their ratio now follows $\gamma=(2+s) / s$. For a lower $T$ one needs only to consider the thermalized degrees of freedom. For a Van der Waals gas holds: $C_{m V}=\frac{1}{2} s R+a p / R T^{2}$.
In general holds:

$$
C_{p}-C_{V}=T\left(\frac{\partial p}{\partial T}\right)_{V} \cdot\left(\frac{\partial V}{\partial T}\right)_{p}=-T\left(\frac{\partial V}{\partial T}\right)_{p}^{2}\left(\frac{\partial p}{\partial V}\right)_{T} \geq 0
$$

Because $(\partial p / \partial V)_{T}$ is always $<0$, the following is always valid: $C_{p} \geq C_{V}$. If the coefficient of expansion is $0, C_{p}=C_{V}$, and also at $T=0 \mathrm{~K}$.

### 8.4 The laws of thermodynamics

The zeroth law states that heat flows from higher to lower temperatures. The first law is the conservation of energy. For a closed system holds: $Q=\Delta U+W$, where $Q$ is the total added heat, $W$ the work done and $\Delta U$ the difference in the internal energy. In differential form this becomes: $\bar{\pi} Q=d U+\pi W$, where $\bar{\pi}$ means that the it is not a differential of a quantity of state. For a quasi-static process holds: $\bar{\pi} W=p d V$. So for a reversible process holds: $\overparen{\pi} Q=d U+p d V$.

For an open (flowing) system the first law is: $Q=\Delta H+W_{\mathrm{i}}+\Delta E_{\mathrm{kin}}+\Delta E_{\mathrm{pot}}$. One can extract an amount of work $W_{\mathrm{t}}$ from the system or add $W_{\mathrm{t}}=-W_{\mathrm{i}}$ to the system.

The second law states: for a closed system there exists an additive quantity $S$, called the entropy, the differential of which has the following property:

$$
d S \geq \frac{\pi Q}{T}
$$

If the only processes occurring are reversible holds: $d S=\pi Q_{\mathrm{rev}} / T$. So, the entropy difference after a reversible process is:

$$
S_{2}-S_{1}=\int_{1}^{2} \frac{\bar{\pi} Q_{\mathrm{rev}}}{T}
$$

So, for a reversible cycle holds: $\oint \frac{\vec{T} Q_{\mathrm{rev}}}{T}=0$.
For an irreversible cycle holds: $\oint \frac{\bar{\pi} Q_{\mathrm{irr}}}{T}<0$.
The third law of thermodynamics is (Nernst):

$$
\lim _{T \rightarrow 0}\left(\frac{\partial S}{\partial X}\right)_{T}=0
$$

From this it can be concluded that the thermal heat capacity $\rightarrow 0$ if $T \rightarrow 0$, so absolute zero temperature cannot be reached by cooling through a finite number of steps.

### 8.5 State functions and Maxwell relations

The quantities of state and their differentials are:

| Internal energy: | $U$ | $d U=T d S-p d V$ |
| :--- | :--- | :--- |
| Enthalpy: | $H=U+p V$ | $d H=T d S+V d p$ |
| Free energy: | $F=U-T S$ | $d F=-S d T-p d V$ |
| Gibbs free enthalpy: | $G=H-T S$ | $d G=-S d T+V d p$ |

From this one can derive Maxwell's relations:

$$
\left(\frac{\partial T}{\partial V}\right)_{S}=-\left(\frac{\partial p}{\partial S}\right)_{V}, \quad\left(\frac{\partial T}{\partial p}\right)_{S}=\left(\frac{\partial V}{\partial S}\right)_{p}, \quad\left(\frac{\partial p}{\partial T}\right)_{V}=\left(\frac{\partial S}{\partial V}\right)_{T}, \quad\left(\frac{\partial V}{\partial T}\right)_{p}=-\left(\frac{\partial S}{\partial p}\right)_{T}
$$

From the total differential and the definitions of $C_{V}$ and $C_{p}$ it can be derived that:

$$
T d S=C_{V} d T+T\left(\frac{\partial p}{\partial T}\right)_{V} d V \text { and } T d S=C_{p} d T-T\left(\frac{\partial V}{\partial T}\right)_{p} d p
$$

For an ideal gas also holds:

$$
S_{m}=C_{V} \ln \left(\frac{T}{T_{0}}\right)+R \ln \left(\frac{V}{V_{0}}\right)+S_{m 0} \text { and } S_{m}=C_{p} \ln \left(\frac{T}{T_{0}}\right)-R \ln \left(\frac{p}{p_{0}}\right)+S_{m 0}^{\prime}
$$

Helmholtz' equations are:

$$
\left(\frac{\partial U}{\partial V}\right)_{T}=T\left(\frac{\partial p}{\partial T}\right)_{V}-p, \quad\left(\frac{\partial H}{\partial p}\right)_{T}=V-T\left(\frac{\partial V}{\partial T}\right)_{p}
$$

for an enlarged surface holds: $\bar{d} W_{\text {rev }}=-\gamma d A$, with $\gamma$ the surface tension. From this follows:

$$
\gamma=\left(\frac{\partial U}{\partial A}\right)_{S}=\left(\frac{\partial F}{\partial A}\right)_{T}
$$

### 8.6 Processes

The efficiency $\eta$ of a process is given by: $\eta=\frac{\text { Work done }}{\text { Heat added }}$
The Cold factor $\xi$ of a cooling down process is given by: $\xi=\frac{\text { Cold delivered }}{\text { Work added }}$

## Reversible adiabatic processes

For adiabatic processes holds: $W=U_{1}-U_{2}$. For reversible adiabatic processes holds Poisson's equation: with $\gamma=C_{p} / C_{V}$ one gets that $p V^{\gamma}=$ constant. Also holds: $T V^{\gamma-1}=\mathrm{constant}$ and $T^{\gamma} p^{1-\gamma}=\mathrm{constant}$. Adiabatics exhibit a greater steepness $p$ - $V$ diagram than isothermics because $\gamma>1$.

## Isobaric processes

Here holds: $H_{2}-H_{1}=\int_{1}^{2} C_{p} d T$. For a reversible isobaric process holds: $H_{2}-H_{1}=Q_{\mathrm{rev}}$.

## The throttle process

This is also called the Joule-Kelvin effect and is an adiabatic expansion of a gas through a porous material or a small opening. Here $H$ is a conserved quantity, and $d S>0$. In general this is accompanied with a change in temperature. The quantity which is important here is the throttle coefficient:

$$
\alpha_{H}=\left(\frac{\partial T}{\partial p}\right)_{H}=\frac{1}{C_{p}}\left[T\left(\frac{\partial V}{\partial T}\right)_{p}-V\right]
$$

The inversion temperature is the temperature where an adiabatically expanding gas keeps the same temperature. If $T>T_{\mathrm{i}}$ the gas heats up, if $T<T_{\mathrm{i}}$ the gas cools down. $T_{\mathrm{i}}=2 T_{\mathrm{B}}$, with for $T_{\mathrm{B}}:[\partial(p V) / \partial p]_{T}=0$. The throttle process is e.g. applied in refridgerators.

## The Carnotprocess

The system undergoes a reversible cycle with 2 isothemics and 2 adiabatics:

1. Isothermic expansion at $T_{1}$. The system absorbs a heat $Q_{1}$ from the reservoir.
2. Adiabatic expansion with a temperature drop to $T_{2}$.
3. Isothermic compression at $T_{2}$, removing $Q_{2}$ from the system.
4. Adiabatic compression to $T_{1}$.

The efficiency for Carnot's process is:

$$
\eta=1-\frac{\left|Q_{2}\right|}{\left|Q_{1}\right|}=1-\frac{T_{2}}{T_{1}}:=\eta_{\mathrm{C}}
$$

The Carnot efficiency $\eta_{\mathrm{C}}$ is the maximal efficiency at which a heat machine can operate. If the process is applied in reverse order and the system performs a work $-W$ the cold factor is given by:

$$
\xi=\frac{\left|Q_{2}\right|}{W}=\frac{\left|Q_{2}\right|}{\left|Q_{1}\right|-\left|Q_{2}\right|}=\frac{T_{2}}{T_{1}-T_{2}}
$$

## The Stirling process

Stirling's cycle exists of 2 isothermics and 2 isochorics. The efficiency in the ideal case is the same as for Carnot's cycle.

### 8.7 Maximal work

Consider a system that changes from state 1 into state 2 , with the temperature and pressure of the surroundings given by $T_{0}$ and $p_{0}$. The maximum work which can be obtained from this change is, when all processes are reversible:

1. Closed system: $W_{\max }=\left(U_{1}-U_{2}\right)-T_{0}\left(S_{1}-S_{2}\right)+p_{0}\left(V_{1}-V_{2}\right)$.
2. Open system: $W_{\max }=\left(H_{1}-H_{2}\right)-T_{0}\left(S_{1}-S_{2}\right)-\Delta E_{\mathrm{kin}}-\Delta E_{\mathrm{pot}}$.

The minimal work needed to attain a certain state is: $W_{\min }=-W_{\max }$.

### 8.8 Phase transitions

Phase transitions are isothermic and isobaric, so $d G=0$. When the phases are indicated by $\alpha, \beta$ and $\gamma$ holds: $G_{m}^{\alpha}=G_{m}^{\beta}$ and

$$
\Delta S_{m}=S_{m}^{\alpha}-S_{m}^{\beta}=\frac{r_{\beta \alpha}}{T_{0}}
$$

where $r_{\beta \alpha}$ is the transition heat of phase $\beta$ to phase $\alpha$ and $T_{0}$ is the transition temperature. The following holds: $r_{\beta \alpha}=r_{\alpha \beta}$ and $r_{\beta \alpha}=r_{\gamma \alpha}-r_{\gamma \beta}$. Further

$$
S_{m}=\left(\frac{\partial G_{m}}{\partial T}\right)_{p}
$$

so $G$ has a twist in the transition point. In a two phase system Clapeyron's equation is valid:

$$
\frac{d p}{d T}=\frac{S_{m}^{\alpha}-S_{m}^{\beta}}{V_{m}^{\alpha}-V_{m}^{\beta}}=\frac{r_{\beta \alpha}}{\left(V_{m}^{\alpha}-V_{m}^{\beta}\right) T}
$$

For an ideal gas one finds for the vapor line at some distance from the critical point:

$$
p=p_{0} \mathrm{e}^{-r_{\beta \alpha / R T}}
$$

There exist also phase transitions with $r_{\beta \alpha}=0$. For those there will occur only a discontinuity in the second derivates of $G_{m}$. These second-order transitions appear at organization phenomena.
A phase-change of the 3 rd order, so with e.g. $\left[\partial^{3} G_{m} / \partial T^{3}\right]_{p}$ non continuous arises e.g. when ferromagnetic iron changes to the paramagnetic state.

### 8.9 Thermodynamic potential

When the number of particles within a system changes this number becomes a third quantity of state. Because addition of matter usually takes place at constant $p$ and $T, G$ is the relevant quantity. If a system exists of more components this becomes:

$$
d G=-S d T+V d p+\sum_{i} \mu_{i} d n_{i}
$$

where $\mu=\left(\frac{\partial G}{\partial n_{i}}\right)_{p, T, n_{j}}$ is called the thermodynamic potential. This is a partial quantity. For $V$ holds:

$$
V=\sum_{i=1}^{c} n_{i}\left(\frac{\partial V}{\partial n_{i}}\right)_{n_{j}, p, T}:=\sum_{i=1}^{c} n_{i} V_{i}
$$

where $V_{i}$ is the partial volume of component $i$. The following holds:

$$
\begin{aligned}
V_{m} & =\sum_{i} x_{i} V_{i} \\
0 & =\sum_{i} x_{i} d V_{i}
\end{aligned}
$$

where $x_{i}=n_{i} / n$ is the molar fraction of component $i$. The molar volume of a mixture of two components can be a concave line in a $V-x_{2}$ diagram: the mixing contracts the volume.

The thermodynamic potentials are not independent in a multiple-phase system. It can be derived that $\sum_{i} n_{i} d \mu_{i}=-S d T+V d p$, this gives at constant $p$ and $T: \sum_{i} x_{i} d \mu_{i}=0$ (Gibbs-Duhmen).
Each component has as much $\mu$ 's as there are phases. The number of free parameters in a system with $c$ components and $p$ different phases is given by $f=c+2-p$.

### 8.10 Ideal mixtures

For a mixture of $n$ components holds (the index ${ }^{0}$ is the value for the pure component):

$$
U_{\text {mixture }}=\sum_{i} n_{i} U_{i}^{0}, H_{\text {mixture }}=\sum_{i} n_{i} H_{i}^{0}, S_{\text {mixture }}=n \sum_{i} x_{i} S_{i}^{0}+\Delta S_{\text {mix }}
$$

where for ideal gases holds: $\Delta S_{\text {mix }}=-n R \sum_{i} x_{i} \ln \left(x_{i}\right)$.
For the thermodynamic potentials holds: $\mu_{i}=\mu_{i}^{0}+R T \ln \left(x_{i}\right)<\mu_{i}^{0}$. A mixture of two liquids is rarely ideal: this is usually only the case for chemically related components or isotopes. In spite of this holds Raoult's law for the vapour pressure holds for many binary mixtures: $p_{i}=x_{i} p_{i}^{0}=y_{i} p$. Here is $x_{i}$ the fraction of the $i$ th component in liquid phase and $y_{i}$ the fraction of the $i$ th component in gas phase.
A solution of one component in another gives rise to an increase in the boiling point $\Delta T_{\mathrm{k}}$ and a decrease of the freezing point $\Delta T_{\mathrm{s}}$. For $x_{2} \ll 1$ holds:

$$
\Delta T_{\mathrm{k}}=\frac{R T_{\mathrm{k}}^{2}}{r_{\beta \alpha}} x_{2}, \Delta T_{\mathrm{s}}=-\frac{R T_{\mathrm{s}}^{2}}{r_{\gamma \beta}} x_{2}
$$

with $r_{\beta \alpha}$ the evaporation heat and $r_{\gamma \beta}<0$ the melting heat. For the osmotic pressure $\Pi$ of a solution holds: $\Pi V_{m 1}^{0}=x_{2} R T$.

### 8.11 Conditions for equilibrium

When a system evolves towards equilibrium the only changes that are possible are those for which holds: $(d S)_{U, V} \geq 0$ or $(d U)_{S, V} \leq 0$ or $(d H)_{S, p} \leq 0$ or $(d F)_{T, V} \leq 0$ or $(d G)_{T, p} \leq 0$. In equilibrium for each component holds: $\mu_{i}^{\alpha}=\mu_{i}^{\beta}=\mu_{i}^{\gamma}$.

### 8.12 Statistical basis for thermodynamics

The number of possibilities $P$ to distribute $N$ particles on $n$ possible energy levels, each with a $g$-fold degeneracy is called the thermodynamic probability and is given by:

$$
P=N!\prod_{i} \frac{g_{i}^{n_{i}}}{n_{i}!}
$$

The most probable distribution, that with the maximum value for $P$, is the equilibrium state. When Stirling's equation, $\ln (n!) \approx n \ln (n)-n$ is used, one finds for a discrete system the Maxwell-Boltzmann distribution. The occupation numbers in equilibrium are then given by:

$$
n_{i}=\frac{N}{Z} g_{i} \exp \left(-\frac{W_{i}}{k T}\right)
$$

The state sum $Z$ is a normalization constant, given by: $Z=\sum_{i} g_{i} \exp \left(-W_{i} / k T\right)$. For an ideal gas holds:

$$
Z=\frac{V(2 \pi m k T)^{3 / 2}}{h^{3}}
$$

The entropy can then be defined as: $S=k \ln (P)$. For a system in thermodynamic equilibrium this becomes:

$$
S=\frac{U}{T}+k N \ln \left(\frac{Z}{N}\right)+k N \approx \frac{U}{T}+k \ln \left(\frac{Z^{N}}{N!}\right)
$$

For an ideal gas, with $U=\frac{3}{2} k T$ then holds: $S=\frac{5}{2} k N+k N \ln \left(\frac{V(2 \pi m k T)^{3 / 2}}{N h^{3}}\right)$

### 8.13 Application to other systems

Thermodynamics can be applied to other systems than gases and liquids. To do this the term $\bar{\pi} W=p d V$ has to be replaced with the correct work term, like $đ W_{\text {rev }}=-F d l$ for the stretching of a wire, $đ W_{\text {rev }}=-\gamma d A$ for the expansion of a soap bubble or $\bar{d} W_{\mathrm{rev}}=-B d M$ for a magnetic system.
A rotating, non-charged black hole has a temparature of $T=\hbar c / 8 \pi k m$. It has an entropy $S=A k c^{3} / 4 \hbar \kappa$ with $A$ the area of its event horizon. For a Schwarzschild black hole $A$ is given by $A=16 \pi \mathrm{~m}^{2}$. Hawkings area theorem states that $d A / d t \geq 0$.

Hence, the lifetime of a black hole $\sim m^{3}$.

## Chapter 9

## Transport phenomena

### 9.1 Mathematical introduction

An important relation is: if $X$ is a quantity of a volume element which travels from position $\vec{r}$ to $\vec{r}+d \vec{r}$ in a time $d t$, the total differential $d X$ is then given by:

$$
d X=\frac{\partial X}{\partial x} d x+\frac{\partial X}{\partial y} d y+\frac{\partial X}{\partial z} d z+\frac{\partial X}{\partial t} d t \Rightarrow \frac{d X}{d t}=\frac{\partial X}{\partial x} v_{x}+\frac{\partial X}{\partial y} v_{y}+\frac{\partial X}{\partial z} v_{z}+\frac{\partial X}{\partial t}
$$

This results in general to: $\frac{d X}{d t}=\frac{\partial X}{\partial t}+(\vec{v} \cdot \nabla) X$.
From this follows that also holds: $\frac{d}{d t} \iiint X d^{3} V=\frac{\partial}{\partial t} \iiint X d^{3} V+\oiint X(\vec{v} \cdot \vec{n}) d^{2} A$
where the volume $V$ is surrounded by surface $A$. Some properties of the $\nabla$ operator are:

$$
\begin{array}{lll}
\operatorname{div}(\phi \vec{v})=\phi \operatorname{div} \vec{v}+\operatorname{grad} \phi \cdot \vec{v} & \operatorname{rot}(\phi \vec{v})=\phi \operatorname{rot} \vec{v}+(\operatorname{grad} \phi) \times \vec{v} & \operatorname{rot} \operatorname{grad} \phi=\overrightarrow{0} \\
\operatorname{div}(\vec{u} \times \vec{v})=\vec{v} \cdot(\operatorname{rot} \vec{u})-\vec{u} \cdot(\operatorname{rot} \vec{v}) & \operatorname{rot} \operatorname{rot} \vec{v}=\operatorname{grad} \operatorname{div} \vec{v}-\nabla^{2} \vec{v} & \operatorname{div} \operatorname{rot} \overrightarrow{\mathrm{v}}=0 \\
\operatorname{div} \operatorname{grad} \phi=\nabla^{2} \phi & \nabla^{2} \vec{v} \equiv\left(\nabla^{2} v_{1}, \nabla^{2} v_{2}, \nabla^{2} v_{3}\right) &
\end{array}
$$

Here, $\vec{v}$ is an arbitrary vector field and $\phi$ an arbitrary scalar field. Some important integral theorems are:
Gauss:

$$
\oiint(\vec{v} \cdot \vec{n}) d^{2} A=\iiint(\operatorname{div} \vec{v}) d^{3} V
$$

Stokes for a scalar field: $\oint\left(\phi \cdot \vec{e}_{\mathrm{t}}\right) d s=\iint(\vec{n} \times \operatorname{grad} \phi) d^{2} A$
Stokes for a vector field: $\oint\left(\vec{v} \cdot \vec{e}_{\mathrm{t}}\right) d s=\iint(\operatorname{rot} \vec{v} \cdot \vec{n}) d^{2} A$
This results in:

$$
\oiint(\operatorname{rot} \vec{v} \cdot \vec{n}) d^{2} A=0
$$

Ostrogradsky:

$$
\begin{aligned}
& \oiint(\vec{n} \times \vec{v}) d^{2} A=\iiint(\operatorname{rot} \vec{v}) d^{3} A \\
& \oiint(\phi \vec{n}) d^{2} A=\iiint(\operatorname{grad} \phi) d^{3} V
\end{aligned}
$$

Here, the orientable surface $\iint d^{2} A$ is limited by the Jordan curve $\oint d s$.

### 9.2 Conservation laws

On a volume work two types of forces:

1. The force $\overrightarrow{f_{0}}$ on each volume element. For gravity holds: $\overrightarrow{f_{0}}=\varrho \vec{g}$.
2. Surface forces working only on the margins: $\vec{t}$. For these holds: $\vec{t}=\vec{n} \mathbf{T}$, where $\mathbf{T}$ is the stress tensor.
$\mathbf{T}$ can be split in a part $p \mathbf{l}$ representing the normal tensions and a part $\mathbf{T}^{\prime}$ representing the shear stresses: $\mathbf{T}=\mathbf{T}^{\prime}+p \mathbf{l}$, where $\mathbf{I}$ is the unit tensor. When viscous aspects can be ignored holds: $\operatorname{div} \mathbf{T}=-\operatorname{grad} p$.
When the flow velocity is $\vec{v}$ at position $\vec{r}$ holds on position $\vec{r}+d \vec{r}$ :

$$
\vec{v}(d \vec{r})=\underbrace{\vec{v}(\vec{r})}_{\text {translation }}+\underbrace{d \vec{r} \cdot(\operatorname{grad} \vec{v})}_{\text {rotation, deformation, dilatation }}
$$

The quantity $\mathbf{L}:=\operatorname{grad} \vec{v}$ can be split in a symmetric part $\mathbf{D}$ and an antisymmetric part $\mathrm{W} . \mathrm{L}=\mathrm{D}+\mathrm{W}$ with

$$
D_{i j}:=\frac{1}{2}\left(\frac{\partial v_{i}}{\partial x_{j}}+\frac{\partial v_{j}}{\partial x_{i}}\right), W_{i j}:=\frac{1}{2}\left(\frac{\partial v_{i}}{\partial x_{j}}-\frac{\partial v_{j}}{\partial x_{i}}\right)
$$

When the rotation or vorticity $\vec{\omega}=\operatorname{rot} \vec{v}$ is introduced holds: $W_{i j}=\frac{1}{2} \varepsilon_{i j k} \omega_{k} . \vec{\omega}$ represents the local rotation velocity: $\overrightarrow{d r} \cdot \mathbf{W}=\frac{1}{2} \omega \times \overrightarrow{d r}$.
For a Newtonian liquid holds: $\mathbf{T}^{\prime}=2 \eta \mathbf{D}$. Here, $\eta$ is the dynamical viscosity. This is related to the shear stress $\tau$ by:

$$
\tau_{i j}=\eta \frac{\partial v_{i}}{\partial x_{j}}
$$

For compressible media can be stated: $\mathbf{T}^{\prime}=\left(\eta^{\prime} \operatorname{div} \vec{v}\right) \mathbf{I}+2 \eta \mathbf{D}$. From equating the thermodynamical and mechanical pressure it follows: $3 \eta^{\prime}+2 \eta=0$. If the viscosity is constant holds: $\operatorname{div}(2 \mathbf{D})=\nabla^{2} \vec{v}+\operatorname{grad} \operatorname{div} \vec{v}$.

The conservation laws for mass, momentum and energy for continuous media can be written in both integral and differential form. They are:

## Integral notation:

1. Conservation of mass: $\frac{\partial}{\partial t} \iiint \varrho d^{3} V+\oiint \varrho(\vec{v} \cdot \vec{n}) d^{2} A=0$
2. Conservation of momentum: $\frac{\partial}{\partial t} \iiint \varrho \stackrel{v}{ } d^{3} V+\oiint \varrho \vec{v}(\vec{v} \cdot \vec{n}) d^{2} A=\iiint f_{0} d^{3} V+\oiint \vec{n} \cdot T d^{2} A$
3. Conservation of energy: $\frac{\partial}{\partial t} \iiint\left(\frac{1}{2} v^{2}+e\right) \varrho d^{3} V+\oiint\left(\frac{1}{2} v^{2}+e\right) \varrho(\vec{v} \cdot \vec{n}) d^{2} A=$

$$
-\oiint(\vec{q} \cdot \vec{n}) d^{2} A+\iiint\left(\vec{v} \cdot \vec{f}_{0}\right) d^{3} V+\oiint(\vec{v} \cdot \vec{n} \mathbf{T}) d^{2} A
$$

## Differential notation:

1. Conservation of mass: $\frac{\partial \varrho}{\partial t}+\operatorname{div} \cdot(\varrho \vec{v})=0$
2. Conservation of momentum: $\varrho \frac{\partial \vec{v}}{\partial t}+(\varrho \vec{v} \cdot \nabla) \vec{v}=\vec{f}_{0}+\operatorname{div} \boldsymbol{T}=\overrightarrow{f_{0}}-\operatorname{grad} p+\operatorname{div} \mathbf{T}^{\prime}$
3. Conservation of energy: $\varrho T \frac{d s}{d t}=\varrho \frac{d e}{d t}-\frac{p}{\varrho} \frac{d \varrho}{d t}=-\operatorname{div} \vec{q}+\mathbf{T}^{\prime}: \mathbf{D}$

Here, $e$ is the internal energy per unit of mass $E / m$ and $s$ is the entropy per unit of mass $S / m \cdot \vec{q}=-\kappa \vec{\nabla} T$ is the heat flow. Further holds:

$$
p=-\frac{\partial E}{\partial V}=-\frac{\partial e}{\partial 1 / \varrho}, \quad T=\frac{\partial E}{\partial S}=\frac{\partial e}{\partial s}
$$

so

$$
C_{V}=\left(\frac{\partial e}{\partial T}\right)_{V} \quad \text { and } \quad C_{p}=\left(\frac{\partial h}{\partial T}\right)_{p}
$$

with $h=H / m$ the enthalpy per unit of mass.

From this one can derive the Navier-Stokes equations for an incompressible, viscous and heat-conducting medium:

$$
\begin{aligned}
\operatorname{div} \vec{v} & =0 \\
\varrho \frac{\partial \vec{v}}{\partial t}+\varrho(\vec{v} \cdot \nabla) \vec{v} & =\varrho \vec{g}-\operatorname{grad} p+\eta \nabla^{2} \vec{v} \\
\varrho C \frac{\partial T}{\partial t}+\varrho C(\vec{v} \cdot \nabla) T & =\kappa \nabla^{2} T+2 \eta \mathbf{D}: \mathbf{D}
\end{aligned}
$$

with $C$ the thermal heat capacity. The force $\vec{F}$ on an object within a flow, when viscous effects are limited to the boundary layer, can be obtained using the momentum law. If a surface $A$ surrounds the object outside the boundary layer holds:

$$
\vec{F}=-\oiint[p \vec{n}+\varrho \vec{v}(\vec{v} \cdot \vec{n})] d^{2} A
$$

### 9.3 Bernoulli's equations

Starting with the momentum equation one can find for a non-viscous medium for stationary flows, with

$$
(\vec{v} \cdot \operatorname{grad}) \vec{v}=\frac{1}{2} \operatorname{grad}\left(v^{2}\right)+(\operatorname{rot} \vec{v}) \times \vec{v}
$$

and the potential equation $\vec{g}=-\operatorname{grad}(g h)$ that:

$$
\frac{1}{2} v^{2}+g h+\int \frac{d p}{\varrho}=\text { constant along a streamline }
$$

For compressible flows holds: $\frac{1}{2} v^{2}+g h+p / \varrho=$ constant along a line of flow. If also holds $\operatorname{rot} \vec{v}=0$ and the entropy is equal on each streamline holds $\frac{1}{2} v^{2}+g h+\int d p / \varrho=$ constant everywhere. For incompressible flows this becomes: $\frac{1}{2} v^{2}+g h+p / \varrho=$ constant everywhere. For ideal gases with constant $C_{p}$ and $C_{V}$ holds, with $\gamma=C_{p} / C_{V}$ :

$$
\frac{1}{2} v^{2}+\frac{\gamma}{\gamma-1} \frac{p}{\varrho}=\frac{1}{2} v^{2}+\frac{c^{2}}{\gamma-1}=\text { constant }
$$

With a velocity potential defined by $\vec{v}=\operatorname{grad} \phi$ holds for instationary flows:

$$
\frac{\partial \phi}{\partial t}+\frac{1}{2} v^{2}+g h+\int \frac{d p}{\varrho}=\text { constant everywhere }
$$

### 9.4 Characterising of flows by dimensionless numbers

The advantage of dimensionless numbers is that they make model experiments possible: one has to make the dimensionless numbers which are important for the specific experiment equal for both model and the real situation. One can also deduce functional equalities without solving the differential equations. Some dimensionless numbers are given by:

$$
\begin{aligned}
& \text { Strouhal: } \quad \mathrm{Sr}=\frac{\omega L}{v} \quad \text { Froude: } \quad \mathrm{Fr}=\frac{v^{2}}{g L} \quad \text { Mach: } \quad \mathrm{Ma}=\frac{v}{c} \\
& \text { Fourier: } \quad \mathrm{Fo}=\frac{a}{\omega L^{2}} \quad \text { Péclet: } \quad \mathrm{Pe}=\frac{v L}{a} \quad \text { Reynolds: } \operatorname{Re}=\frac{v L}{\nu} \\
& \text { Prandtl: } \quad \operatorname{Pr}=\frac{\nu}{a} \quad \text { Nusselt: } \quad \mathrm{Nu}=\frac{L \alpha}{\kappa} \quad \text { Eckert: } \quad \mathrm{Ec}=\frac{v^{2}}{c \Delta T}
\end{aligned}
$$

Here, $\nu=\eta / \varrho$ is the kinematic viscosity, $c$ is the speed of sound and $L$ is a characteristic length of the system. $\alpha$ follows from the equation for heat transport $\kappa \partial_{y} T=\alpha \Delta T$ and $a=\kappa / \varrho c$ is the thermal diffusion coefficient.
These numbers can be interpreted as follows:

- Re: (stationary inertial forces)/(viscous forces)
- Sr: (non-stationary inertial forces)/(stationary inertial forces)
- Fr: (stationary inertial forces)/(gravity)
- Fo: (heat conductance)/(non-stationary change in enthalpy)
- Pe: (convective heat transport)/(heat conductance)
- Ec: (viscous dissipation)/(convective heat transport)
- Ma: (velocity)/(speed of sound): objects moving faster than approximately $\mathrm{Ma}=0,8$ produce shockwaves which propagate with an angle $\theta$ with the velocity of the object. For this angle holds $\mathrm{Ma}=$ $1 / \arctan (\theta)$.
- Pr and Nu are related to specific materials.

Now, the dimensionless Navier-Stokes equation becomes, with $x^{\prime}=x / L, \vec{v}^{\prime}=\vec{v} / V$, grad $=L$ grad, $\nabla^{\prime 2}=$ $L^{2} \nabla^{2}$ and $t^{\prime}=t \omega$ :

$$
\operatorname{Sr} \frac{\partial \vec{v}^{\prime}}{\partial t^{\prime}}+\left(\vec{v}^{\prime} \cdot \nabla^{\prime}\right) \vec{v}^{\prime}=-\operatorname{grad}^{\prime} p+\frac{\vec{g}}{\operatorname{Fr}}+\frac{\nabla^{\prime 2} \vec{v}^{\prime}}{\operatorname{Re}}
$$

### 9.5 Tube flows

For tube flows holds: they are laminar if $\mathrm{Re}<2300$ with dimension of length the diameter of the tube, and turbulent if Re is larger. For an incompressible laminar flow through a straight, circular tube holds for the velocity profile:

$$
v(r)=-\frac{1}{4 \eta} \frac{d p}{d x}\left(R^{2}-r^{2}\right)
$$

For the volume flow holds: $\Phi_{V}=\int_{0}^{R} v(r) 2 \pi r d r=-\frac{\pi}{8 \eta} \frac{d p}{d x} R^{4}$
The entrance length $L_{\mathrm{e}}$ is given by:

1. $500<\operatorname{Re}_{D}<2300: L_{\mathrm{e}} / 2 R=0.056 \operatorname{Re}_{D}$
2. $\mathrm{Re}>2300: L_{\mathrm{e}} / 2 R \approx 50$

For gas transport at low pressures (Knudsen-gas) holds: $\Phi_{V}=\frac{4 R^{3} \alpha \sqrt{\pi}}{3} \frac{d p}{d x}$
For flows at a small Re holds: $\nabla p=\eta \nabla^{2} \vec{v}$ and $\operatorname{div} \vec{v}=0$. For the total force on a sphere with radius $R$ in a flow then holds: $F=6 \pi \eta R v$. For large Re holds for the force on a surface $A$ : $F=\frac{1}{2} C_{W} A \varrho v^{2}$.

### 9.6 Potential theory

The circulation $\Gamma$ is defined as: $\Gamma=\oint\left(\vec{v} \cdot \overrightarrow{e_{\mathrm{t}}}\right) d s=\iint(\operatorname{rot} \vec{v}) \cdot \vec{n} d^{2} A=\iint(\vec{\omega} \cdot \vec{n}) d^{2} A$
For non viscous media, if $p=p(\varrho)$ and all forces are conservative, Kelvin's theorem can be derived:

$$
\frac{d \Gamma}{d t}=0
$$

For rotationless flows a velocity potential $\vec{v}=\operatorname{grad} \phi$ can be introduced. In the incompressible case follows from conservation of mass $\nabla^{2} \phi=0$. For a 2 -dimensional flow a flow function $\psi(x, y)$ can be defined: with $\Phi_{A B}$ the amount of liquid flowing through a curve $s$ between the points A and B:

$$
\Phi_{A B}=\int_{A}^{B}(\vec{v} \cdot \vec{n}) d s=\int_{A}^{B}\left(v_{x} d y-v_{y} d x\right)
$$

and the definitions $v_{x}=\partial \psi / \partial y, v_{y}=-\partial \psi / \partial x$ holds: $\Phi_{A B}=\psi(B)-\psi(A)$. In general holds:

$$
\frac{\partial^{2} \psi}{\partial x^{2}}+\frac{\partial^{2} \psi}{\partial y^{2}}=-\omega_{z}
$$

In polar coordinates holds:

$$
v_{r}=\frac{1}{r} \frac{\partial \psi}{\partial \theta}=\frac{\partial \phi}{\partial r}, v_{\theta}=-\frac{\partial \psi}{\partial r}=\frac{1}{r} \frac{\partial \phi}{\partial \theta}
$$

For source flows with power $Q$ in $(x, y)=(0,0)$ holds: $\phi=\frac{Q}{2 \pi} \ln (r)$ so that $v_{r}=Q / 2 \pi r, v_{\theta}=0$.
For a dipole of strength $Q$ in $x=a$ and strength $-Q$ in $x=-a$ follows from superposition: $\phi=-Q a x / 2 \pi r^{2}$ where $Q a$ is the dipole strength. For a vortex holds: $\phi=\Gamma \theta / 2 \pi$.

If an object is surrounded by an uniform main flow with $\vec{v}=v \vec{e}_{x}$ and such a large Re that viscous effects are limited to the boundary layer holds: $F_{x}=0$ and $F_{y}=-\varrho \Gamma v$. The statement that $F_{x}=0$ is d'Alembert's paradox and originates from the neglection of viscous effects. The lift $F_{y}$ is also created by $\eta$ because $\Gamma \neq 0$ due to viscous effects. Henxe rotating bodies also create a force perpendicular to their direction of motion: the Magnus effect.

### 9.7 Boundary layers

### 9.7.1 Flow boundary layers

If for the thickness of the boundary layer holds: $\delta \ll L$ holds: $\delta \approx L / \sqrt{\operatorname{Re}}$. With $v_{\infty}$ the velocity of the main flow it follows for the velocity $v_{y} \perp$ the surface: $v_{y} L \approx \delta v_{\infty}$. Blasius' equation for the boundary layer is, with $v_{y} / v_{\infty}=f(y / \delta): 2 f^{\prime \prime \prime}+f f^{\prime \prime}=0$ with boundary conditions $f(0)=f^{\prime}(0)=0, f^{\prime}(\infty)=1$. From this follows: $C_{W}=0.664 \mathrm{Re}_{x}^{-1 / 2}$.
The momentum theorem of Von Karman for the boundary layer is: $\frac{d}{d x}\left(\vartheta v^{2}\right)+\delta^{*} v \frac{d v}{d x}=\frac{\tau_{0}}{\varrho}$
where the displacement thickness $\delta^{*} v$ and the momentum thickness $\vartheta v^{2}$ are given by:

$$
\vartheta v^{2}=\int_{0}^{\infty}\left(v-v_{x}\right) v_{x} d y, \quad \delta^{*} v=\int_{0}^{\infty}\left(v-v_{x}\right) d y \text { and } \tau_{0}=-\left.\eta \frac{\partial v_{x}}{\partial y}\right|_{y=0}
$$

The boundary layer is released from the surface if $\left(\frac{\partial v_{x}}{\partial y}\right)_{y=0}=0$. This is equivalent with $\frac{d p}{d x}=\frac{12 \eta v_{\infty}}{\delta^{2}}$.

### 9.7.2 Temperature boundary layers

If the thickness of the temperature boundary layer $\delta_{T} \ll L$ holds: 1 . If $\operatorname{Pr} \leq 1: \delta / \delta_{T} \approx \sqrt{\operatorname{Pr}}$.
2. If $\operatorname{Pr} \gg 1: \delta / \delta_{T} \approx \sqrt[3]{\operatorname{Pr}}$.

### 9.8 Heat conductance

For non-stationairy heat conductance in one dimension without flow holds:

$$
\frac{\partial T}{\partial t}=\frac{\kappa}{\varrho c} \frac{\partial^{2} T}{\partial x^{2}}+\Phi
$$

where $\Phi$ is a source term. If $\Phi=0$ the solutions for harmonic oscillations at $x=0$ are:

$$
\frac{T-T_{\infty}}{T_{\max }-T_{\infty}}=\exp \left(-\frac{x}{D}\right) \cos \left(\omega t-\frac{x}{D}\right)
$$

with $D=\sqrt{2 \kappa / \omega \varrho c}$. At $x=\pi D$ the temperature variation is in anti-phase with the surface. The onedimensional solution at $\Phi=0$ is

$$
T(x, t)=\frac{1}{2 \sqrt{\pi a t}} \exp \left(-\frac{x^{2}}{4 a t}\right)
$$

This is mathematical equivalent to the diffusion problem:

$$
\frac{\partial n}{\partial t}=D \nabla^{2} n+P-A
$$

where $P$ is the production of and $A$ the discharge of particles. The flow density $J=-D \nabla n$.

### 9.9 Turbulence

The time scale of turbulent velocity variations $\tau_{\mathrm{t}}$ is of the order of: $\tau_{\mathrm{t}}=\tau \sqrt{\mathrm{Re}} / \mathrm{Ma}^{2}$ with $\tau$ the molecular time scale. For the velocity of the particles holds: $v(t)=\langle v\rangle+v^{\prime}(t)$ with $\left\langle v^{\prime}(t)\right\rangle=0$. The Navier-Stokes equation now becomes:

$$
\frac{\partial\langle\vec{v}\rangle}{\partial t}+(\langle\vec{v}\rangle \cdot \nabla)\langle\vec{v}\rangle=-\frac{\nabla\langle p\rangle}{\varrho}+\nu \nabla^{2}\langle\vec{v}\rangle+\frac{\operatorname{div} \mathrm{S}_{R}}{\varrho}
$$

where $\mathbf{S}_{R i j}=-\varrho\left\langle v_{i} v_{j}\right\rangle$ is the turbulent stress tensor. Boussinesq's assumption is: $\tau_{i j}=-\varrho\left\langle v_{i}^{\prime} v_{j}^{\prime}\right\rangle$. It is stated that, analogous to Newtonian media: $\mathbf{S}_{R}=2 \varrho \nu_{t}\langle\mathbf{D}\rangle$. Near a boundary holds: $\nu_{t}=0$, far away of a boundary holds: $\nu_{t} \approx \nu \operatorname{Re}$.

### 9.10 Self organization

For a (semi) two-dimensional flow holds: $\frac{d \omega}{d t}=\frac{\partial \omega}{\partial t}+J(\omega, \psi)=\nu \nabla^{2} \omega$
With $J(\omega, \psi)$ the Jacobian. So if $\nu=0, \omega$ is conserved. Further, the kinetic energy/ $m A$ and the enstrofy $V$ are conserved: with $\vec{v}=\nabla \times(\vec{k} \psi)$

$$
E \sim(\nabla \psi)^{2} \sim \int_{0}^{\infty} \mathcal{E}(k, t) d k=\text { constant }, \quad V \sim\left(\nabla^{2} \psi\right)^{2} \sim \int_{0}^{\infty} k^{2} \mathcal{E}(k, t) d k=\text { constant }
$$

From this follows that in a two-dimensional flow the energy flux goes towards large values of $k$ : larger structures become larger at the expanse of smaller ones. In three-dimensional flows the situation is just the opposite.

## Chapter 10

## Quantum physics

### 10.1 Introduction to quantum physics

### 10.1.1 Black body radiation

Planck's law for the energy distribution for the radiation of a black body is:

$$
w(f)=\frac{8 \pi h f^{3}}{c^{3}} \frac{1}{\mathrm{e}^{h f / k T}-1}, \quad w(\lambda)=\frac{8 \pi h c}{\lambda^{5}} \frac{1}{\mathrm{e}^{h c / \lambda k T}-1}
$$

Stefan-Boltzmann's law for the total power density can be derived from this: $P=A \sigma T^{4}$. Wien's law for the maximum can also be derived from this: $T \lambda_{\max }=k_{\mathrm{W}}$.

### 10.1.2 The Compton effect

For the wavelength of scattered light, if light is considered to exist of particles, can be derived:

$$
\lambda^{\prime}=\lambda+\frac{h}{m c}(1-\cos \theta)=\lambda+\lambda_{\mathrm{C}}(1-\cos \theta)
$$

### 10.1.3 Electron diffraction

Diffraction of electrons at a crystal can be explained by assuming that particles have a wave character with wavelength $\lambda=h / p$. This wavelength is called the Broglie-wavelength.

### 10.2 Wave functions

The wave character of particles is described by a wavefunction $\psi$. This wavefunction can be described in normal or momentum space. Both definitions are each others Fourier transform:

$$
\Phi(k, t)=\frac{1}{\sqrt{h}} \int \Psi(x, t) \mathrm{e}^{-i k x} d x \text { and } \Psi(x, t)=\frac{1}{\sqrt{h}} \int \Phi(k, t) \mathrm{e}^{i k x} d k
$$

These waves define a particle with group velocity $v_{\mathrm{g}}=p / m$ and energy $E=\hbar \omega$.
The wavefunction can be interpreted as a measure for the probability $P$ to find a particle somewhere (Born): $d P=|\psi|^{2} d^{3} V$. The expectation value $\langle f\rangle$ of a quantity $f$ of a system is given by:

$$
\langle f(t)\rangle=\iiint \Psi^{*} f \Psi d^{3} V, \quad\left\langle f_{p}(t)\right\rangle=\iiint \Phi^{*} f \Phi d^{3} V_{p}
$$

This is also written as $\langle f(t)\rangle=\langle\Phi| f|\Phi\rangle$. The normalizing condition for wavefunctions follows from this: $\langle\Phi \mid \Phi\rangle=\langle\Psi \mid \Psi\rangle=1$.

### 10.3 Operators in quantum physics

In quantum mechanics, classical quantities are translated into operators. These operators are hermitian because their eigenvalues must be real:

$$
\int \psi_{1}^{*} A \psi_{2} d^{3} V=\int \psi_{2}\left(A \psi_{1}\right)^{*} d^{3} V
$$

When $u_{n}$ is the eigenfunction of the eigenvalue equation $A \Psi=a \Psi$ for eigenvalue $a_{n}, \Psi$ can be expanded into a basis of eigenfunctions: $\Psi=\sum_{n} c_{n} u_{n}$. If this basis is taken orthonormal, then follows for the coefficients: $c_{n}=\left\langle u_{n} \mid \Psi\right\rangle$. If the system is in a state described by $\Psi$, the chance to find eigenvalue $a_{n}$ when measuring $A$ is given by $\left|c_{n}\right|^{2}$ in the discrete part of the spectrum and $\left|c_{n}\right|^{2} d a$ in the continuous part of the spectrum between $a$ and $a+d a$. The matrix element $A_{i j}$ is given by: $A_{i j}=\left\langle u_{i}\right| A\left|u_{j}\right\rangle$. Because $(A B)_{i j}=\left\langle u_{i}\right| A B\left|u_{j}\right\rangle=$ $\left\langle u_{i}\right| A \sum_{n}\left|u_{n}\right\rangle\left\langle u_{n}\right| B\left|u_{j}\right\rangle$ holds: $\sum_{n}\left|u_{n}\right\rangle\left\langle u_{n}\right|=1$.
The time-dependence of an operator is given by (Heisenberg):

$$
\frac{d A}{d t}=\frac{\partial A}{\partial t}+\frac{[A, H]}{i \hbar}
$$

with $[A, B] \equiv A B-B A$ the commutator of $A$ and $B$. For hermitian operators the commutator is always complex. If $[A, B]=0$, the operators $A$ and $B$ have a common set of eigenfunctions. By applying this to $p_{x}$ and $x$ follows (Ehrenfest): $m d^{2}\langle x\rangle_{t} / d t^{2}=-\langle d U(x) / d x\rangle$.
The first order approximation $\langle F(x)\rangle_{t} \approx F(\langle x\rangle)$, with $F=-d U / d x$ represents the classical equation.
Before the addition of quantummechanical operators which are a product of other operators, they should be made symmetrical: a classical product $A B$ becomes $\frac{1}{2}(A B+B A)$.

### 10.4 The uncertainty principle

If the uncertainty $\Delta A$ in $A$ is defined as: $\left.(\Delta A)^{2}=\langle\psi| A_{\text {op }}-\left.\langle A\rangle\right|^{2} \psi\right\rangle=\left\langle A^{2}\right\rangle-\langle A\rangle^{2}$ it follows:

$$
\left.\Delta A \cdot \Delta B \geq \frac{1}{2}|\langle\psi|[A, B]| \psi\right\rangle \mid
$$

From this follows: $\Delta E \cdot \Delta t \geq \frac{1}{2} \hbar$, and because $\left[x, p_{x}\right]=i \hbar$ holds: $\Delta p_{x} \cdot \Delta x \geq \frac{1}{2} \hbar$, and $\Delta L_{x} \cdot \Delta L_{y} \geq \frac{1}{2} \hbar L_{z}$.

### 10.5 The Schrödinger equation

The momentum operator is given by: $p_{\mathrm{op}}=-i \hbar \nabla$. The position operator is: $x_{\mathrm{op}}=i \hbar \nabla_{p}$. The energy operator is given by: $E_{\mathrm{op}}=i \hbar \partial / \partial t$. The Hamiltonian of a particle with mass $m$, potential energy $U$ and total energy $E$ is given by: $H=p^{2} / 2 m+U$. From $H \psi=E \psi$ then follows the Schrödinger equation:

$$
-\frac{\hbar^{2}}{2 m} \nabla^{2} \psi+U \psi=E \psi=i \hbar \frac{\partial \psi}{\partial t}
$$

The linear combination of the solutions of this equation give the general solution. In one dimension it is:

$$
\psi(x, t)=\left(\sum+\int d E\right) c(E) u_{E}(x) \exp \left(-\frac{i E t}{\hbar}\right)
$$

The current density $J$ is given by: $J=\frac{\hbar}{2 i m}\left(\psi^{*} \nabla \psi-\psi \nabla \psi^{*}\right)$
The following conservation law holds: $\frac{\partial P(x, t)}{\partial t}=-\nabla J(x, t)$

### 10.6 Parity

The parity operator in one dimension is given by $\mathcal{P} \psi(x)=\psi(-x)$. If the wavefunction is split in even and odd functions, it can be expanded into eigenfunctions of $\mathcal{P}$ :

$$
\psi(x)=\underbrace{\frac{1}{2}(\psi(x)+\psi(-x))}_{\text {even: } \psi^{+}}+\underbrace{\frac{1}{2}(\psi(x)-\psi(-x))}_{\text {odd: } \psi^{-}}
$$

$[\mathcal{P}, H]=0$. The functions $\psi^{+}=\frac{1}{2}(1+\mathcal{P}) \psi(x, t)$ and $\psi^{-}=\frac{1}{2}(1-\mathcal{P}) \psi(x, t)$ both satisfy the Schrödinger equation. Hence, parity is a conserved quantity.

### 10.7 The tunnel effect

The wavefunction of a particle in an $\infty$ high potential step from $x=0$ to $x=a$ is given by $\psi(x)=$ $a^{-1 / 2} \sin (k x)$. The energylevels are given by $E_{n}=n^{2} h^{2} / 8 a^{2} m$.
If the wavefunction with energy $W$ meets a potential well of $W_{0}>W$ the wavefunction will, unlike the classical case, be non-zero within the potential well. If 1, 2 and 3 are the areas in front, within and behind the potential well, holds:

$$
\psi_{1}=A \mathrm{e}^{i k x}+B \mathrm{e}^{-i k x}, \quad \psi_{2}=C \mathrm{e}^{i k^{\prime} x}+D \mathrm{e}^{-i k^{\prime} x}, \quad \psi_{3}=A^{\prime} \mathrm{e}^{i k x}
$$

with $k^{\prime 2}=2 m\left(W-W_{0}\right) / \hbar^{2}$ and $k^{2}=2 m W$. Using the boundary conditions requiring continuity: $\psi=$ continuous and $\partial \psi / \partial x=$ continuous at $x=0$ and $x=a$ gives $B, C$ and $D$ and $A^{\prime}$ expressed in $A$. The amplitude $T$ of the transmitted wave is defined by $T=\left|A^{\prime}\right|^{2} /|A|^{2}$. If $W>W_{0}$ and $2 a=n \lambda^{\prime}=2 \pi n / k^{\prime}$ holds: $T=1$.

### 10.8 The harmonic oscillator

For a harmonic oscillator holds: $U=\frac{1}{2} b x^{2}$ and $\omega_{0}^{2}=b / m$. The Hamiltonian $H$ is then given by:

$$
H=\frac{p^{2}}{2 m}+\frac{1}{2} m \omega^{2} x^{2}=\frac{1}{2} \hbar \omega+\omega A^{\dagger} A
$$

with

$$
A=\sqrt{\frac{1}{2} m \omega} x+\frac{i p}{\sqrt{2 m \omega}} \text { and } A^{\dagger}=\sqrt{\frac{1}{2} m \omega} x-\frac{i p}{\sqrt{2 m \omega}}
$$

$A \neq A^{\dagger}$ is non hermitian. $\left[A, A^{\dagger}\right]=\hbar$ and $[A, H]=\hbar \omega A$. $A$ is a so called raising ladder operator, $A^{\dagger}$ a lowering ladder operator. $H A u_{E}=(E-\hbar \omega) A u_{E}$. There is an eigenfunction $u_{0}$ for which holds: $A u_{0}=0$. The energy in this ground state is $\frac{1}{2} \hbar \omega$ : the zero point energy. For the normalized eigenfunctions follows:

$$
u_{n}=\frac{1}{\sqrt{n!}}\left(\frac{A^{\dagger}}{\sqrt{\hbar}}\right)^{n} u_{0} \text { with } u_{0}=\sqrt[4]{\frac{m \omega}{\pi \hbar}} \exp \left(-\frac{m \omega x^{2}}{2 \hbar}\right)
$$

with $E_{n}=\left(\frac{1}{2}+n\right) \hbar \omega$.

### 10.9 Angular momentum

For the angular momentum operators $L$ holds: $\left[L_{z}, L^{2}\right]=\left[L_{z}, H\right]=\left[L^{2}, H\right]=0$. However, cyclically holds: $\left[L_{x}, L_{y}\right]=i \hbar L_{z}$. Not all components of $L$ can be known at the same time with arbitrary accuracy. For $L_{z}$ holds:

$$
L_{z}=-i \hbar \frac{\partial}{\partial \varphi}=-i \hbar\left(x \frac{\partial}{\partial y}-y \frac{\partial}{\partial x}\right)
$$

The ladder operators $L_{ \pm}$are defined by: $L_{ \pm}=L_{x} \pm i L_{y}$. Now holds: $L^{2}=L_{+} L_{-}+L_{z}^{2}-\hbar L_{z}$. Further,

$$
L_{ \pm}=\hbar \mathrm{e}^{ \pm i \varphi}\left( \pm \frac{\partial}{\partial \theta}+i \cot (\theta) \frac{\partial}{\partial \varphi}\right)
$$

From $\left[L_{+}, L_{z}\right]=-\hbar L_{+}$follows: $L_{z}\left(L_{+} Y_{l m}\right)=(m+1) \hbar\left(L_{+} Y_{l m}\right)$.
From $\left[L_{-}, L_{z}\right]=\hbar L_{-}$follows: $L_{z}\left(L_{-} Y_{l m}\right)=(m-1) \hbar\left(L_{-} Y_{l m}\right)$.
From $\left[L^{2}, L_{ \pm}\right]=0$ follows: $L^{2}\left(L_{ \pm} Y_{l m}\right)=l(l+1) \hbar^{2}\left(L_{ \pm} Y_{l m}\right)$.
Because $L_{x}$ and $L_{y}$ are hermitian (this implies $L_{ \pm}^{\dagger}=L_{\mp}$ ) and $\left|L_{ \pm} Y_{l m}\right|^{2}>0$ follows: $l(l+1)-m^{2}-m \geq$ $0 \Rightarrow-l \leq m \leq l$. Further follows that $l$ has to be integral or half-integral. Half-odd integral values give no unique solution $\psi$ and are therefore dismissed.

### 10.10 Spin

For the spin operators are defined by their commutation relations: $\left[S_{x}, S_{y}\right]=i \hbar S_{z}$. Because the spin operators do not act in the physical space $(x, y, z)$ the uniqueness of the wavefunction is not a criterium here: also half odd-integer values are allowed for the spin. Because $[L, S]=0$ spin and angular momentum operators do not have a common set of eigenfunctions. The spin operators are given by $\overrightarrow{\vec{S}}=\frac{1}{2} \hbar \overrightarrow{\vec{\sigma}}$, with

$$
\vec{\sigma}_{x}=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right) \quad, \quad \overrightarrow{\vec{\sigma}}_{y}=\left(\begin{array}{cc}
0 & -i \\
i & 0
\end{array}\right) \quad, \quad \overrightarrow{\vec{\sigma}}_{z}=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)
$$

The eigenstates of $S_{z}$ are called spinors: $\chi=\alpha_{+} \chi_{+}+\alpha_{-} \chi_{-}$, where $\chi_{+}=(1,0)$ represents the state with spin up $\left(S_{z}=\frac{1}{2} \hbar\right)$ and $\chi_{-}=(0,1)$ represents the state with spin down $\left(S_{z}=-\frac{1}{2} \hbar\right)$. Then the probability to find spin up after a measurement is given by $\left|\alpha_{+}\right|^{2}$ and the chance to find spin down is given by $\left|\alpha_{-}\right|^{2}$. Of course holds $\left|\alpha_{+}\right|^{2}+\left|\alpha_{-}\right|^{2}=1$.
The electron will have an intrinsic magnetic dipole moment $\vec{M}$ due to its spin, given by $\vec{M}=-e g_{S} \vec{S} / 2 m$, with $g_{S}=2(1+\alpha / 2 \pi+\cdots)$ the gyromagnetic ratio. In the presence of an external magnetic field this gives a potential energy $U=-\vec{M} \cdot \vec{B}$. The Schrödinger equation then becomes (because $\partial \chi / \partial x_{i} \equiv 0$ ):

$$
i \hbar \frac{\partial \chi(t)}{\partial t}=\frac{e g_{S} \hbar}{4 m} \vec{\sigma} \cdot \vec{B} \chi(t)
$$

with $\vec{\sigma}=\left(\overrightarrow{\vec{\sigma}}_{x}, \overrightarrow{\vec{\sigma}}_{y}, \overrightarrow{\vec{\sigma}}_{z}\right)$. If $\vec{B}=B \vec{e}_{z}$ there are two eigenvalues for this problem: $\chi_{ \pm}$for $E= \pm e g_{S} \hbar B / 4 m=$ $\pm \hbar \omega$. So the general solution is given by $\chi=\left(a \mathrm{e}^{-i \omega t}, b \mathrm{e}^{i \omega t}\right)$. From this can be derived: $\left\langle S_{x}\right\rangle=\frac{1}{2} \hbar \cos (2 \omega t)$ and $\left\langle S_{y}\right\rangle=\frac{1}{2} \hbar \sin (2 \omega t)$. Thus the spin precesses about the $z$-axis with frequency $2 \omega$. This causes the normal Zeeman splitting of spectral lines.
The potential operator for two particles with spin $\pm \frac{1}{2} \hbar$ is given by:

$$
V(r)=V_{1}(r)+\frac{1}{\hbar^{2}}\left(\vec{S}_{1} \cdot \overrightarrow{S_{2}}\right) V_{2}(r)=V_{1}(r)+\frac{1}{2} V_{2}(r)\left[S(S+1)-\frac{3}{2}\right]
$$

This makes it possible for two states to exist: $S=1$ (triplet) or $S=0$ (Singlet).

### 10.11 The Dirac formalism

If the operators for $p$ and $E$ are substituted in the relativistic equation $E^{2}=m_{0}^{2} c^{4}+p^{2} c^{2}$, the Klein-Gordon equation is found:

$$
\left(\nabla^{2}-\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}-\frac{m_{0}^{2} c^{2}}{\hbar^{2}}\right) \psi(\vec{x}, t)=0
$$

The operator $\square-m_{0}^{2} c^{2} / \hbar^{2}$ can be separated:

$$
\nabla^{2}-\frac{1}{c^{2}} \frac{\partial^{2}}{\partial t^{2}}-\frac{m_{0}^{2} c^{2}}{\hbar^{2}}=\left\{\gamma_{\lambda} \frac{\partial}{\partial x_{\lambda}}-\frac{m_{0}^{2} c^{2}}{\hbar^{2}}\right\}\left\{\gamma_{\mu} \frac{\partial}{\partial x_{\mu}}+\frac{m_{0}^{2} c^{2}}{\hbar^{2}}\right\}
$$

where the Dirac matrices $\gamma$ are given by: $\gamma_{\lambda} \gamma_{\mu}+\gamma_{\mu} \gamma_{\lambda}=2 \delta_{\lambda \mu}$. From this it can be derived that the $\gamma$ are hermitian $4 \times 4$ matrices given by:

$$
\gamma_{k}=\left(\begin{array}{cc}
0 & -i \sigma_{k} \\
i \sigma_{k} & 0
\end{array}\right) \quad, \quad \gamma_{4}=\left(\begin{array}{cc}
I & 0 \\
0 & -I
\end{array}\right)
$$

With this, the Dirac equation becomes:

$$
\left(\gamma_{\lambda} \frac{\partial}{\partial x_{\lambda}}+\frac{m_{0}^{2} c^{2}}{\hbar^{2}}\right) \psi(\vec{x}, t)=0
$$

where $\psi(x)=\left(\psi_{1}(x), \psi_{2}(x), \psi_{3}(x), \psi_{4}(x)\right)$ is a spinor.

### 10.12 Atomic physics

### 10.12.1 Solutions

The solutions of the Schrödinger equation in spherical coordinates if the potential energy is a function of $r$ alone can be written as: $\psi(r, \theta, \varphi)=R_{n l}(r) Y_{l, m_{l}}(\theta, \varphi) \chi_{m_{s}}$, with

$$
Y_{l m}=\frac{C_{l m}}{\sqrt{2 \pi}} P_{l}^{m}(\cos \theta) \mathrm{e}^{i m \varphi}
$$

For an atom or ion with one electron holds: $R_{l m}(\rho)=C_{l m} \mathrm{e}^{-\rho / 2} \rho^{l} L_{n-l-1}^{2 l+1}(\rho)$
with $\rho=2 r Z / n a_{0}$ with $a_{0}=\varepsilon_{0} h^{2} / \pi m_{\mathrm{e}} e^{2}$. The $L_{i}^{j}$ are the associated Laguere functions and the $P_{l}^{m}$ are the associated Legendre polynomials:

$$
P_{l}^{|m|}(x)=\left(1-x^{2}\right)^{m / 2} \frac{d^{|m|}}{d x^{|m|}}\left[\left(x^{2}-1\right)^{l}\right] \quad, \quad L_{n}^{m}(x)=\frac{(-1)^{m} n!}{(n-m)!} \mathrm{e}^{-x} x^{-m} \frac{d^{n-m}}{d x^{n-m}}\left(\mathrm{e}^{-x} x^{n}\right)
$$

The parity of these solutions is $(-1)^{l}$. The functions are $2 \sum_{l=0}^{n-1}(2 l+1)=2 n^{2}$-folded degenerated.

### 10.12.2 Eigenvalue equations

The eigenvalue equations for an atom or ion with with one electron are:

| Equation | Eigenvalue | Range |
| :--- | :--- | :--- |
| $H_{\mathrm{op}} \psi=E \psi$ | $E_{n}=\mu e^{4} Z^{2} / 8 \varepsilon_{0}^{2} h^{2} n^{2}$ | $n \geq 1$ |
| $L_{z \mathrm{op}} Y_{l m}=L_{z} Y_{l m}$ | $L_{z}=m_{l} \hbar$ | $-l \leq m_{l} \leq l$ |
| $L_{\mathrm{op}}^{2} Y_{l m}=L^{2} Y_{l m}$ | $L^{2}=l(l+1) \hbar^{2}$ | $l<n$ |
| $S_{z \mathrm{op}} \chi=S_{z} \chi$ | $S_{z}=m_{s} \hbar$ | $m_{s}= \pm \frac{1}{2}$ |
| $S_{\mathrm{op}}^{2} \chi=S^{2} \chi$ | $S^{2}=s(s+1) \hbar^{2}$ | $s=\frac{1}{2}$ |

### 10.12.3 Spin-orbit interaction

The total momentum is given by $\vec{J}=\vec{L}+\vec{M}$. The total magnetic dipole moment of an electron is then $\vec{M}=\vec{M}_{L}+\vec{M}_{S}=-\left(e / 2 m_{\mathrm{e}}\right)\left(\vec{L}+g_{S} \vec{S}\right)$ where $g_{S}=2.0023$ is the gyromagnetic ratio of the electron. Further holds: $J^{2}=L^{2}+S^{2}+2 \vec{L} \cdot \vec{S}=L^{2}+S^{2}+2 L_{z} S_{z}+L_{+} S_{-}+L_{-} S_{+} . J$ has quantum numbers $j$ with possible values $j=l \pm \frac{1}{2}$, with $2 j+1$ possible $z$-components ( $m_{J} \in\{-j, \ldots, 0, . ., j\}$ ). If the interaction energy between $S$ and $L$ is small it can be stated that: $E=E_{n}+E_{S L}=E_{n}+a \vec{S} \cdot \vec{L}$. It can then be derived that:

$$
a=\frac{\left|E_{n}\right| Z^{2} \alpha^{2}}{\hbar^{2} n l(l+1)\left(l+\frac{1}{2}\right)}
$$

After a relativistic correction this becomes:

$$
E=E_{n}+\frac{\left|E_{n}\right| Z^{2} \alpha^{2}}{n}\left(\frac{3}{4 n}-\frac{1}{j+\frac{1}{2}}\right)
$$

The fine structure in atomic spectra arises from this. With $g_{S}=2$ follows for the average magnetic moment: $\vec{M}_{\mathrm{av}}=-\left(e / 2 m_{\mathrm{e}}\right) g \hbar \vec{J}$, where $g$ is the Landé-factor:

$$
g=1+\frac{\vec{S} \cdot \vec{J}}{J^{2}}=1+\frac{j(j+1)+s(s+1)-l(l+1)}{2 j(j+1)}
$$

For atoms with more than one electron the following limiting situations occur:

1. $L-S$ coupling: for small atoms the electrostatic interaction is dominant and the state can be characterized by $L, S, J, m_{J} . J \in\{|L-S|, \ldots, L+S-1, L+S\}$ and $m_{J} \in\{-J, \ldots, J-1, J\}$. The spectroscopic notation for this interaction is: ${ }^{2 S+1} L_{J} .2 S+1$ is the multiplicity of a multiplet.
2. $j-j$ coupling: for larger atoms the electrostatic interaction is smaller than the $L_{i} \cdot s_{i}$ interaction of an electron. The state is characterized by $j_{i} \ldots j_{n}, J, m_{J}$ where only the $j_{i}$ of the not completely filled subshells are to be taken into account.

The energy difference for larger atoms when placed in a magnetic field is: $\Delta E=g \mu_{\mathrm{B}} m_{J} B$ where $g$ is the Landé factor. For a transition between two singlet states the line splits in 3 parts, for $\Delta m_{J}=-1,0+1$. This results in the normal Zeeman effect. At higher $S$ the line splits up in more parts: the anomalous Zeeman effect.

Interaction with the spin of the nucleus gives the hyperfine structure.

### 10.12.4 Selection rules

For the dipole transition matrix elements follows: $\left.p_{0} \sim\left|\left\langle l_{2} m_{2}\right| \vec{E} \cdot \vec{r}\right| l_{1} m_{1}\right\rangle \mid$. Conservation of angular momentum demands that for the transition of an electron holds that $\Delta l= \pm 1$.

For an atom where $L-S$ coupling is dominant further holds: $\Delta S=0$ (but not strict), $\Delta L=0, \pm 1, \Delta J=$ $0, \pm 1$ except for $J=0 \rightarrow J=0$ transitions, $\Delta m_{J}=0, \pm 1$, but $\Delta m_{J}=0$ is forbidden if $\Delta J=0$.

For an atom where $j-j$ coupling is dominant further holds: for the jumping electron holds, except $\Delta l= \pm 1$, also: $\Delta j=0, \pm 1$, and for all other electrons: $\Delta j=0$. For the total atom holds: $\Delta J=0, \pm 1$ but no $J=0 \rightarrow J=0$ transitions and $\Delta m_{J}=0, \pm 1$, but $\Delta m_{J}=0$ is forbidden if $\Delta J=0$.

### 10.13 Interaction with electromagnetic fields

The Hamiltonian of an electron in an electromagnetic field is given by:

$$
H=\frac{1}{2 \mu}(\vec{p}+e \vec{A})^{2}-e V=-\frac{\hbar^{2}}{2 \mu} \nabla^{2}+\frac{e}{2 \mu} \vec{B} \cdot \vec{L}+\frac{e^{2}}{2 \mu} A^{2}-e V
$$

where $\mu$ is the reduced mass of the system. The term $\sim A^{2}$ can usually be neglected, except for very strong fields or macroscopic motions. For $\vec{B}=B \vec{e}_{z}$ it is given by $e^{2} B^{2}\left(x^{2}+y^{2}\right) / 8 \mu$.

When a gauge transformation $\overrightarrow{A^{\prime}}=\vec{A}-\nabla f, V^{\prime}=V+\partial f / \partial t$ is applied to the potentials the wavefunction is also transformed according to $\psi^{\prime}=\psi \mathrm{e}^{i q e f} / \hbar$ with $q e$ the charge of the particle. Because $f=f(x, t)$, this is called a local gauge transformation, in contrast with a global gauge transformation which can always be applied.

### 10.14 Perturbation theory

### 10.14.1 Time-independent perturbation theory

To solve the equation $\left(H_{0}+\lambda H_{1}\right) \psi_{n}=E_{n} \psi_{n}$ one has to find the eigenfunctions of $H=H_{0}+\lambda H_{1}$. Suppose that $\phi_{n}$ is a complete set of eigenfunctions of the non-perturbed Hamiltonian $H_{0}: H_{0} \phi_{n}=E_{n}^{0} \phi_{n}$. Because $\phi_{n}$ is a complete set holds:

$$
\psi_{n}=N(\lambda)\left\{\phi_{n}+\sum_{k \neq n} c_{n k}(\lambda) \phi_{k}\right\}
$$

When $c_{n k}$ and $E_{n}$ are being expanded into $\lambda: c_{n k}=\lambda c_{n k}^{(1)}+\lambda^{2} c_{n k}^{(2)}+\cdots$

$$
E_{n}=E_{n}^{0}+\lambda E_{n}^{(1)}+\lambda^{2} E_{n}^{(2)}+\cdots
$$

and this is put into the Schrödinger equation the result is: $E_{n}^{(1)}=\left\langle\phi_{n}\right| H_{1}\left|\phi_{n}\right\rangle$ and $c_{n m}^{(1)}=\frac{\left\langle\phi_{m}\right| H_{1}\left|\phi_{n}\right\rangle}{E_{n}^{0}-E_{m}^{0}}$ if $m \neq n$. The second-order correction of the energy is then given by: $E_{n}^{(2)}=\sum_{k \neq n} \frac{\left.\left|\left\langle\phi_{k}\right| H_{1}\right| \phi_{n}\right\rangle\left.\right|^{2}}{E_{n}^{0}-E_{k}^{0}}$. So to first order holds: $\psi_{n}=\phi_{n}+\sum_{k \neq n} \frac{\left\langle\phi_{k}\right| \lambda H_{1}\left|\phi_{n}\right\rangle}{E_{n}^{0}-E_{k}^{0}} \phi_{k}$.

In case the levels are degenerated the above does not hold. In that case an orthonormal set eigenfunctions $\phi_{n i}$ is chosen for each level $n$, so that $\left\langle\phi_{m i} \mid \phi_{n j}\right\rangle=\delta_{m n} \delta_{i j}$. Now $\psi$ is expanded as:

$$
\psi_{n}=N(\lambda)\left\{\sum_{i} \alpha_{i} \phi_{n i}+\lambda \sum_{k \neq n} c_{n k}^{(1)} \sum_{i} \beta_{i} \phi_{k i}+\cdots\right\}
$$

$E_{n i}=E_{n i}^{0}+\lambda E_{n i}^{(1)}$ is approximated by $E_{n i}^{0}:=E_{n}^{0}$. Substitution in the Schrödinger equation and taking dot product with $\phi_{n i}$ gives: $\sum_{i} \alpha_{i}\left\langle\phi_{n j}\right| H_{1}\left|\phi_{n i}\right\rangle=E_{n}^{(1)} \alpha_{j}$. Normalization requires that $\sum_{i}\left|\alpha_{i}\right|^{2}=1$.

### 10.14.2 Time-dependent perturbation theory

From the Schrödinger equation $i \hbar \frac{\partial \psi(t)}{\partial t}=\left(H_{0}+\lambda V(t)\right) \psi(t)$
and the expansion $\psi(t)=\sum_{n} c_{n}(t) \exp \left(\frac{-i E_{n}^{0} t}{\hbar}\right) \phi_{n}$ with $c_{n}(t)=\delta_{n k}+\lambda c_{n}^{(1)}(t)+\cdots$
follows: $c_{n}^{(1)}(t)=\frac{\lambda}{i \hbar} \int_{0}^{t}\left\langle\phi_{n}\right| V\left(t^{\prime}\right)\left|\phi_{k}\right\rangle \exp \left(\frac{i\left(E_{n}^{0}-E_{k}^{0}\right) t^{\prime}}{\hbar}\right) d t^{\prime}$

### 10.15 N-particle systems

### 10.15.1 General

Identical particles are indistinguishable. For the total wavefunction of a system of identical indistinguishable particles holds:

1. Particles with a half-odd integer spin (Fermions): $\psi_{\text {total }}$ must be antisymmetric w.r.t. interchange of the coordinates (spatial and spin) of each pair of particles. The Pauli principle results from this: two Fermions cannot exist in an identical state because then $\psi_{\text {total }}=0$.
2. Particles with an integer spin (Bosons): $\psi_{\text {total }}$ must be symmetric w.r.t. interchange of the coordinates (spatial and spin) of each pair of particles.

For a system of two electrons there are 2 possibilities for the spatial wavefunction. When $a$ and $b$ are the quantum numbers of electron 1 and 2 holds:

$$
\psi_{\mathrm{S}}(1,2)=\psi_{a}(1) \psi_{b}(2)+\psi_{a}(2) \psi_{b}(1), \psi_{\mathrm{A}}(1,2)=\psi_{a}(1) \psi_{b}(2)-\psi_{a}(2) \psi_{b}(1)
$$

Because the particles do not approach each other closely the repulsion energy at $\psi_{\mathrm{A}}$ in this state is smaller. The following spin wavefunctions are possible:

$$
\begin{aligned}
\chi_{\mathrm{A}} & =\frac{1}{2} \sqrt{2}\left[\chi_{+}(1) \chi_{-}(2)-\chi_{+}(2) \chi_{-}(1)\right] \\
m_{s} & =0 \\
\chi_{\mathrm{S}} & = \begin{cases}\chi_{+}(1) \chi_{+}(2) & m_{s}=+1 \\
\frac{1}{2} \sqrt{2}\left[\chi_{+}(1) \chi_{-}(2)+\chi_{+}(2) \chi_{-}(1)\right] & m_{s}=0 \\
\chi_{-}(1) \chi_{-}(2) & m_{s}=-1\end{cases}
\end{aligned}
$$

Because the total wavefunction must be antisymmetric it follows: $\psi_{\text {total }}=\psi_{\mathrm{S}} \chi_{\mathrm{A}}$ or $\psi_{\text {total }}=\psi_{\mathrm{A}} \chi_{\mathrm{S}}$.

For $N$ particles the symmetric spatial function is given by:

$$
\psi_{\mathrm{S}}(1, \ldots, N)=\sum \psi(\text { all permutations of } 1 . . N)
$$

The antisymmetric wavefunction is given by the determinant $\psi_{\mathrm{A}}(1, \ldots, N)=\frac{1}{\sqrt{N!}}\left|u_{E_{i}}(j)\right|$

### 10.15.2 Molecules

The wavefunctions of atom $a$ and $b$ are $\phi_{a}$ and $\phi_{b}$. If the 2 atoms approach each other there are two possibilities: the total wavefunction approaches the bonding function with lower total energy $\psi_{\mathrm{B}}=\frac{1}{2} \sqrt{2}\left(\phi_{a}+\phi_{b}\right)$ or approaches the anti-bonding function with higher energy $\psi_{\mathrm{AB}}=\frac{1}{2} \sqrt{2}\left(\phi_{a}-\phi_{b}\right)$. If a molecular-orbital is symmetric w.r.t. the connecting axis, like a combination of two s-orbitals it is called a $\sigma$-orbital, otherwise a $\pi$-orbital, like the combination of two p-orbitals along two axes.
The energy of a system is: $E=\frac{\langle\psi| H|\psi\rangle}{\langle\psi \mid \psi\rangle}$.
The energy calculated with this method is always higher than the real energy if $\psi$ is only an approximation for the solutions of $H \psi=E \psi$. Also, if there are more functions to be chosen, the function which gives the lowest energy is the best approximation. Applying this to the function $\psi=\sum c_{i} \phi_{i}$ one finds: $\left(H_{i j}-E S_{i j}\right) c_{i}=0$. This equation has only solutions if the secular determinant $\left|H_{i j}-E S_{i j}\right|=0$. Here, $H_{i j}=\left\langle\phi_{i}\right| H\left|\phi_{j}\right\rangle$ and $S_{i j}=\left\langle\phi_{i} \mid \phi_{j}\right\rangle . \alpha_{i}:=H_{i i}$ is the Coulomb integral and $\beta_{i j}:=H_{i j}$ the exchange integral. $S_{i i}=1$ and $S_{i j}$ is the overlap integral.
The first approximation in the molecular-orbital theory is to place both electrons of a chemical bond in the bonding orbital: $\psi(1,2)=\psi_{\mathrm{B}}(1) \psi_{\mathrm{B}}(2)$. This results in a large electron density between the nuclei and therefore a repulsion. A better approximation is: $\psi(1,2)=C_{1} \psi_{\mathrm{B}}(1) \psi_{\mathrm{B}}(2)+C_{2} \psi_{\mathrm{AB}}(1) \psi_{\mathrm{AB}}(2)$, with $C_{1}=1$ and $C_{2} \approx 0.6$.
In some atoms, such as $C$, it is energetical more suitable to form orbitals which are a linear combination of the $\mathrm{s}, \mathrm{p}$ and d states. There are three ways of hybridization in C :

1. SP-hybridization: $\psi_{\mathrm{sp}}=\frac{1}{2} \sqrt{2}\left(\psi_{2 \mathrm{~s}} \pm \psi_{2 \mathrm{p}_{z}}\right)$. There are 2 hybrid orbitals which are placed on one line under $180^{\circ}$. Further the $2 \mathrm{p}_{x}$ and $2 \mathrm{p}_{y}$ orbitals remain.
2. $\mathrm{SP}^{2}$ hybridization: $\psi_{\mathrm{sp}^{2}}=\psi_{2 \mathrm{~s}} / \sqrt{3}+c_{1} \psi_{2 \mathrm{p}_{z}}+c_{2} \psi_{2 \mathrm{p}_{y}}$, where $\left(c_{1}, c_{2}\right) \in\{(\sqrt{2 / 3}, 0),(-1 / \sqrt{6}, 1 / \sqrt{2})$ $,(-1 / \sqrt{6},-1 / \sqrt{2})\}$. The $3 \mathrm{SP}^{2}$ orbitals lay in one plane, with symmetry axes which are at an angle of $120^{\circ}$.
3. $\mathrm{SP}^{3}$ hybridization: $\psi_{\mathrm{sp}^{3}}=\frac{1}{2}\left(\psi_{2 \mathrm{~s}} \pm \psi_{2 \mathrm{p}_{z}} \pm \psi_{2 \mathrm{p}_{y}} \pm \psi_{2 \mathrm{p}_{x}}\right)$. The $4 \mathrm{SP}^{3}$ orbitals form a tetraheder with the symmetry axes at an angle of $109^{\circ} 28^{\prime}$.

### 10.16 Quantum statistics

If a system exists in a state in which one has not the disposal of the maximal amount of information about the system, it can be described by a density matrix $\rho$. If the probability that the system is in state $\psi_{i}$ is given by $a_{i}$, one can write for the expectation value $a$ of $A:\langle a\rangle=\sum_{i} r_{i}\left\langle\psi_{i}\right| A\left|\psi_{i}\right\rangle$.
If $\psi$ is expanded into an orthonormal basis $\left\{\phi_{k}\right\}$ as: $\psi^{(i)}=\sum_{k} c_{k}^{(i)} \phi_{k}$, holds:

$$
\langle A\rangle=\sum_{k}(A \rho)_{k k}=\operatorname{Tr}(A \rho)
$$

where $\rho_{l k}=c_{k}^{*} c_{l} . \rho$ is hermitian, with $\operatorname{Tr}(\rho)=1$. Further holds $\rho=\sum r_{i}\left|\psi_{i}\right\rangle\left\langle\psi_{i}\right|$. The probability to find eigenvalue $a_{n}$ when measuring $A$ is given by $\rho_{n n}$ if one uses a basis of eigenvectors of $A$ for $\left\{\phi_{k}\right\}$. For the time-dependence holds (in the Schrödinger image operators are not explicitly time-dependent):

$$
i \hbar \frac{d \rho}{d t}=[H, \rho]
$$

For a macroscopic system in equilibrium holds $[H, \rho]=0$. If all quantumstates with the same energy are equally probable: $P_{i}=P\left(E_{i}\right)$, one can obtain the distribution:

$$
P_{n}(E)=\rho_{n n}=\frac{\mathrm{e}^{-E_{n} / k T}}{Z} \text { with the state sum } Z=\sum_{n} \mathrm{e}^{-E_{n} / k T}
$$

The thermodynamic quantities are related to these definitions as follows: $F=-k T \ln (Z), U=\langle H\rangle=$ $\sum_{n} p_{n} E_{n}=-\frac{\partial}{\partial k T} \ln (Z), S=-k \sum_{n} P_{n} \ln \left(P_{n}\right)$. For a mixed state of $M$ orthonormal quantum states with probability $1 / M$ follows: $S=k \ln (M)$.
The distribution function for the internal states for a system in thermal equilibrium is the most probable function. This function can be found by taking the maximum of the function which gives the number of states with Stirling's equation: $\ln (n!) \approx n \ln (n)-n$, and the conditions $\sum_{k} n_{k}=N$ and $\sum_{k} n_{k} W_{k}=W$. For identical, indistinguishable particles which obey the Pauli exclusion principle the possible number of states is given by:

$$
P=\prod_{k} \frac{g_{k}!}{n_{k}!\left(g_{k}-n_{k}\right)!}
$$

This results in the Fermi-Dirac statistics. For indistinguishable particles which do not obey the exclusion principle the possible number of states is given by:

$$
P=N!\prod_{k} \frac{g_{k}^{n_{k}}}{n_{k}!}
$$

This results in the Bose-Einstein statistics. So the distribution functions which explain how particles are distributed over the different one-particle states $k$ which are each $g_{k}$-fold degenerate depend on the spin of the particles. They are given by:

1. Fermi-Dirac statistics: integer spin. $n_{k} \in\{0,1\}, n_{k}=\frac{N}{Z_{g}} \frac{g_{k}}{\exp \left(\left(E_{k}-\mu\right) / k T\right)+1}$
with $\ln \left(Z_{\mathrm{g}}\right)=\sum g_{k} \ln \left[1+\exp \left(\left(E_{i}-\mu\right) / k T\right)\right]$.
2. Bose-Einstein statistics: half odd-integer spin. $n_{k} \in I N, n_{k}=\frac{N}{Z_{g}} \frac{g_{k}}{\exp \left(\left(E_{k}-\mu\right) / k T\right)-1}$ with $\ln \left(Z_{\mathrm{g}}\right)=-\sum g_{k} \ln \left[1-\exp \left(\left(E_{i}-\mu\right) / k T\right)\right]$.
Here, $Z_{\mathrm{g}}$ is the large-canonical state sum and $\mu$ the chemical potential. It is found by demanding $\sum n_{k}=N$, and for it holds: $\lim _{T \rightarrow 0} \mu=E_{\mathrm{F}}$, the Fermi-energy. $N$ is the total number of particles. The Maxwell-Boltzmann distribution can be derived from this in the limit $E_{k}-\mu \gg k T$ :

$$
n_{k}=\frac{N}{Z} \exp \left(-\frac{E_{k}}{k T}\right) \text { with } Z=\sum_{k} g_{k} \exp \left(-\frac{E_{k}}{k T}\right)
$$

With the Fermi-energy, the Fermi-Dirac and Bose-Einstein statistics can be written as:

1. Fermi-Dirac statistics: $n_{k}=\frac{g_{k}}{\exp \left(\left(E_{k}-E_{\mathrm{F}}\right) / k T\right)+1}$.
2. Bose-Einstein statistics: $n_{k}=\frac{g_{k}}{\exp \left(\left(E_{k}-E_{\mathrm{F}}\right) / k T\right)-1}$.

## Chapter 11

## Plasma physics

### 11.1 Introduction

The degree of ionization $\alpha$ of a plasma is defined by: $\alpha=\frac{n_{\mathrm{e}}}{n_{\mathrm{e}}+n_{0}}$
where $n_{\mathrm{e}}$ is the electron density and $n_{0}$ the density of the neutrals. If a plasma contains also negative charged ions $\alpha$ is not well defined.

The probability that a test particle collides with another is given by $d P=n \sigma d x$ where $\sigma$ is the cross section. The collision frequency $\nu_{\mathrm{c}}=1 / \tau_{\mathrm{c}}=n \sigma v$. The mean free path is given by $\lambda_{\mathrm{v}}=1 / n \sigma$. The rate coefficient $K$ is defined by $K=\langle\sigma v\rangle$. The number of collisions per unit of time and volume between particles of kind 1 and 2 is given by $n_{1} n_{2}\langle\sigma v\rangle=K n_{1} n_{2}$.
The potential of an electron is given by:

$$
V(r)=\frac{-e}{4 \pi \varepsilon_{0} r} \exp \left(-\frac{r}{\lambda_{\mathrm{D}}}\right) \quad \text { with } \quad \lambda_{\mathrm{D}}=\sqrt{\frac{\varepsilon_{0} k T_{\mathrm{e}} T_{\mathrm{i}}}{e^{2}\left(n_{\mathrm{e}} T_{\mathrm{i}}+n_{\mathrm{i}} T_{\mathrm{e}}\right)}} \approx \sqrt{\frac{\varepsilon_{0} k T_{\mathrm{e}}}{n_{\mathrm{e}} e^{2}}}
$$

because charge is shielded in a plasma. Here, $\lambda_{\mathrm{D}}$ is the Debye length. For distances $<\lambda_{\mathrm{D}}$ the plasma cannot be assumed to be quasi-neutral. Deviations of charge neutrality by thermic motion are compensated by oscillations with frequency

$$
\omega_{\mathrm{pe}}=\sqrt{\frac{n_{\mathrm{e}} e^{2}}{m_{\mathrm{e}} \varepsilon_{0}}}
$$

The distance of closest approximation when two equal charged particles collide for a deviation of $\pi / 2$ is $2 b_{0}=e^{2} /\left(4 \pi \varepsilon_{0} \frac{1}{2} m v^{2}\right)$. A "neat" plasma is defined as a plasma for which holds: $b_{0}<n_{\mathrm{e}}^{-1 / 3} \ll \lambda_{\mathrm{D}} \ll L_{\mathrm{p}}$. Here $L_{\mathrm{p}}:=\left|n_{\mathrm{e}} / \nabla n_{\mathrm{e}}\right|$ is the gradient length of the plasma.

### 11.2 Transport

Relaxation times are defined as $\tau=1 / \nu_{\mathrm{c}}$. Starting with $\sigma_{\mathrm{m}}=4 \pi b_{0}^{2} \ln \left(\Lambda_{\mathrm{C}}\right)$ and with $\frac{1}{2} m v^{2}=k T$ it can be found that:

$$
\tau_{\mathrm{m}}=\frac{4 \pi \varepsilon_{0}^{2} m^{2} v^{3}}{n e^{4} \ln \left(\Lambda_{\mathrm{C}}\right)}=\frac{8 \sqrt{2} \pi \varepsilon_{0}^{2} \sqrt{m}(k T)^{3 / 2}}{n e^{4} \ln \left(\Lambda_{\mathrm{C}}\right)}
$$

For momentum transfer between electrons and ions holds for a Maxwellian velocity distribution:

$$
\tau_{\mathrm{ee}}=\frac{6 \pi \sqrt{3} \varepsilon_{0}^{2} \sqrt{m_{\mathrm{e}}}\left(k T_{\mathrm{e}}\right)^{3 / 2}}{n_{\mathrm{e}} e^{4} \ln \left(\Lambda_{\mathrm{C}}\right)} \approx \tau_{\mathrm{ei}}, \quad \tau_{\mathrm{ii}}=\frac{6 \pi \sqrt{3} \varepsilon_{0}^{2} \sqrt{m_{\mathrm{i}}}\left(k T_{\mathrm{i}}\right)^{3 / 2}}{n_{\mathrm{i}} e^{4} \ln \left(\Lambda_{\mathrm{C}}\right)}
$$

The energy relaxation times for identical particles are equal to the momentum relaxation times. Because for e-i collisions the energy transfer is only $\sim 2 m_{\mathrm{e}} / m_{\mathrm{i}}$ this is a slow process. Approximately holds: $\tau_{\text {ee }}: \tau_{\text {ei }}$ : $\tau_{\text {ie }}: \tau_{\text {ie }}^{E}=1: 1: \sqrt{m_{\mathrm{i}} / m_{\mathrm{e}}}: m_{\mathrm{i}} / m_{\mathrm{e}}$.
The relaxation for e-o interaction is much more complicated. For $T>10 \mathrm{eV}$ holds approximately: $\sigma_{\mathrm{eo}}=$ $10^{-17} v_{\mathrm{e}}^{-2 / 5}$, for lower energies this can be a factor 10 lower.
The resistivity $\eta=E / J$ of a plasma is given by:

$$
\eta=\frac{n_{\mathrm{e}} e^{2}}{m_{\mathrm{e}} \nu_{\mathrm{ei}}}=\frac{e^{2} \sqrt{m_{\mathrm{e}}} \ln \left(\Lambda_{\mathrm{C}}\right)}{6 \pi \sqrt{3} \varepsilon_{0}^{2}\left(k T_{\mathrm{e}}\right)^{3 / 2}}
$$

The diffusion coefficient $D$ is defined by means of the flux $\Gamma$ by $\vec{\Gamma}=n \vec{v}_{\text {diff }}=-D \nabla n$. The equation of continuity is $\partial_{t} n+\nabla\left(n v_{\mathrm{diff}}\right)=0 \Rightarrow \partial_{t} n=D \nabla^{2} n$. One finds that $D=\frac{1}{3} \lambda_{\mathrm{v}} v$. A rough estimate gives $\tau_{\mathrm{D}}=L_{\mathrm{p}} / D=L_{\mathrm{p}}^{2} \tau_{\mathrm{c}} / \lambda_{\mathrm{v}}^{2}$. For magnetized plasma's $\lambda_{\mathrm{v}}$ must be replaced with the cyclotron radius. In electrical fields also holds $\vec{J}=n e \mu \vec{E}=e\left(n_{\mathrm{e}} \mu_{\mathrm{e}}+n_{\mathrm{i}} \mu_{\mathrm{i}}\right) \vec{E}$ with $\mu=e / m \nu_{\mathrm{c}}$ the mobility of the particles. The Einstein ratio is:

$$
\frac{D}{\mu}=\frac{k T}{e}
$$

Because a plasma is electrically neutral electrons and ions are strongly coupled and they don't diffuse independent. The coefficient of ambipolar diffusion $D_{\mathrm{amb}}$ is defined by $\vec{\Gamma}=\vec{\Gamma}_{\mathrm{i}}=\vec{\Gamma}_{\mathrm{e}}=-D_{\mathrm{amb}} \nabla n_{\mathrm{e}, \mathrm{i}}$. From this follows that

$$
D_{\mathrm{amb}}=\frac{k T_{\mathrm{e}} / e-k T_{\mathrm{i}} / e}{1 / \mu_{\mathrm{e}}-1 / \mu_{\mathrm{i}}} \approx \frac{k T_{\mathrm{e}} \mu_{\mathrm{i}}}{e}
$$

In an external magnetic field $B_{0}$ particles will move in spiral orbits with cyclotron radius $\rho=m v / e B_{0}$ and with cyclotron frequency $\Omega=B_{0} e / m$. The helical orbit is perturbed by collisions. A plasma is called magnetized if $\lambda_{\mathrm{v}}>\rho_{\mathrm{e}, \mathrm{i}}$. So the electrons are magnetized if

$$
\frac{\rho_{\mathrm{e}}}{\lambda_{\mathrm{ee}}}=\frac{\sqrt{m_{\mathrm{e}}} e^{3} n_{\mathrm{e}} \ln \left(\Lambda_{\mathrm{C}}\right)}{6 \pi \sqrt{3} \varepsilon_{0}^{2}\left(k T_{\mathrm{e}}\right)^{3 / 2} B_{0}}<1
$$

Magnetization of only the electrons is sufficient to confine the plasma reasonable because they are coupled to the ions by charge neutrality. In case of magnetic confinement holds: $\nabla p=\vec{J} \times \vec{B}$. Combined with the two stationary Maxwell equations for the $B$-field these form the ideal magneto-hydrodynamic equations. For a uniform $B$-field holds: $p=n k T=B^{2} / 2 \mu_{0}$.

If both magnetic and electric fields are present electrons and ions will move in the same direction. If $\vec{E}=$ $E_{r} \vec{e}_{r}+E_{z} \vec{e}_{z}$ and $\vec{B}=B_{z} \vec{e}_{z}$ the $\vec{E} \times \vec{B}$ drift results in a velocity $\vec{u}=(\vec{E} \times \vec{B}) / B^{2}$ and the velocity in the $r, \varphi$ plane is $\dot{r}(r, \varphi, t)=\vec{u}+\dot{\vec{\rho}}(t)$.

### 11.3 Elastic collisions

### 11.3.1 General

The scattering angle of a particle in interaction with another particle, as shown in the figure at the right is:

$$
\chi=\pi-2 b \int_{r_{a}}^{\infty} \frac{d r}{r^{2} \sqrt{1-\frac{b^{2}}{r^{2}}-\frac{W(r)}{E_{0}}}}
$$

Particles with an impact parameter between $b$ and $b+d b$,
 moving through a ring with $d \sigma=2 \pi b d b$ leave the scattering area at a solid angle $d \Omega=2 \pi \sin (\chi) d \chi$. The differential cross section is then defined as:

$$
I(\Omega)=\left|\frac{d \sigma}{d \Omega}\right|=\frac{b}{\sin (\chi)} \frac{\partial b}{\partial \chi}
$$

For a potential energy $W(r)=k r^{-n}$ follows: $I(\Omega, v) \sim v^{-4 / n}$.
For low energies, $\mathcal{O}(1 \mathrm{eV}), \sigma$ has a Ramsauer minimum. It arises from the interference of matter waves behind the object. $I(\Omega)$ for angles $0<\chi<\lambda / 4$ is larger than the classical value.

### 11.3.2 The Coulomb interaction

For the Coulomb interaction holds: $2 b_{0}=q_{1} q_{2} / 2 \pi \varepsilon_{0} m v_{0}^{2}$, so $W(r)=2 b_{0} / r$. This gives $b=b_{0} \cot \left(\frac{1}{2} \chi\right)$ and

$$
I\left(\Omega=\frac{b}{\sin (\chi)} \frac{\partial b}{\partial \chi}=\frac{b_{0}^{2}}{4 \sin ^{2}\left(\frac{1}{2} \chi\right)}\right.
$$

Because the influence of a particle vanishes at $r=\lambda_{\mathrm{D}}$ holds: $\sigma=\pi\left(\lambda_{\mathrm{D}}^{2}-b_{0}^{2}\right)$. Because $d p=d(m v)=$ $m v_{0}(1-\cos \chi)$ a cross section related to momentum transfer $\sigma_{\mathrm{m}}$ is given by:

$$
\sigma_{\mathrm{m}}=\int(1-\cos \chi) I(\Omega) d \Omega=4 \pi b_{0}^{2} \ln \left(\frac{1}{\sin \left(\frac{1}{2} \chi_{\min }\right)}\right)=4 \pi b_{0}^{2} \ln \left(\frac{\lambda_{\mathrm{D}}}{b_{0}}\right):=4 \pi b_{0}^{2} \ln \left(\Lambda_{\mathrm{C}}\right) \sim \frac{\ln \left(v^{4}\right)}{v^{4}}
$$

where $\ln \left(\Lambda_{\mathrm{C}}\right)$ is the Coulomb-logarithm. For this quantity holds: $\Lambda_{\mathrm{C}}=\lambda_{\mathrm{D}} / b_{0}=9 n\left(\lambda_{\mathrm{D}}\right)$.

### 11.3.3 The induced dipole interaction

The induced dipole interaction, with $\vec{p}=\alpha \vec{E}$, gives a potential $V$ and an energy $W$ in a dipole field given by:

$$
V(r)=\frac{\vec{p} \cdot \vec{e}_{r}}{4 \pi \varepsilon_{0} r^{2}}, W(r)=-\frac{|e| p}{8 \pi \varepsilon_{0} r^{2}}=-\frac{\alpha e^{2}}{2\left(4 \pi \varepsilon_{0}\right)^{2} r^{4}}
$$

with $b_{a}=\sqrt[4]{\frac{2 e^{2} \alpha}{\left(4 \pi \varepsilon_{0}\right)^{2} \frac{1}{2} m v_{0}^{2}}}$ holds: $\chi=\pi-2 b \int_{r_{a}}^{\infty} \frac{d r}{r^{2} \sqrt{1-\frac{b^{2}}{r^{2}}+\frac{b_{a}^{4}}{4 r^{4}}}}$
If $b \geq b_{a}$ the charge would hit the atom. Repulsing nuclear forces prevent this to happen. If the scattering angle is a lot times $2 \pi$ it is called capture. The cross section for capture $\sigma_{\text {orb }}=\pi b_{a}^{2}$ is called the Langevin limit, and is a lowest estimate for the total cross section.

### 11.3.4 The centre of mass system

If collisions of two particles with masses $m_{1}$ and $m_{2}$ which scatter in the centre of mass system by an angle $\chi$ are compared with the scattering under an angle $\theta$ in the laboratory system holds:

$$
\tan (\theta)=\frac{m_{2} \sin (\chi)}{m_{1}+m_{2} \cos (\chi)}
$$

The energy loss $\Delta E$ of the incoming particle is given by:

$$
\frac{\Delta E}{E}=\frac{\frac{1}{2} m_{2} v_{2}^{2}}{\frac{1}{2} m_{1} v_{1}^{2}}=\frac{2 m_{1} m_{2}}{\left(m_{1}+m_{2}\right)^{2}}(1-\cos (\chi))
$$

### 11.3.5 Scattering of light

Scattering of light by free electrons is called Thomson scattering. The scattering is free from collective effects if $k \lambda_{\mathrm{D}} \ll 1$. The cross section $\sigma=6.65 \cdot 10^{-29} \mathrm{~m}^{2}$ and

$$
\frac{\Delta f}{f}=\frac{2 v}{c} \sin \left(\frac{1}{2} \chi\right)
$$

This gives for the scattered energy $E_{\text {scat }} \sim n \lambda_{0}^{4} /\left(\lambda^{2}-\lambda_{0}^{2}\right)^{2}$ with $n$ the density. If $\lambda \gg \lambda_{0}$ it is called Rayleigh scattering. Thomson sccattering is a limit of Compton scattering, which is given by $\lambda^{\prime}-\lambda=\lambda_{\mathrm{C}}(1-\cos \chi)$ with $\lambda_{\mathrm{C}}=h / m c$ and cannot be used any more if relativistic effects become important.

### 11.4 Thermodynamic equilibrium and reversibility

Planck's radiation law and the Maxwellian velocity distribution hold for a plasma in equilibrium:

$$
\rho(\nu, T) d \nu=\frac{8 \pi h \nu^{3}}{c^{3}} \frac{1}{\exp (h \nu / k T)-1} d \nu \quad, \quad N(E, T) d E=\frac{2 \pi n}{(\pi k T)^{3 / 2}} \sqrt{E} \exp \left(-\frac{E}{k T}\right) d E
$$

"Detailed balancing" means that the number of reactions in one direction equals the number of reactions in the opposite direction because both processes have equal probability if one corrects for the used phase space. For the reaction

$$
\sum_{\text {forward }} \mathrm{X}_{\text {forward }} \rightleftarrows \sum_{\text {back }} \mathrm{X}_{\text {back }}
$$

holds in a plasma in equilibrium microscopic reversibility:

$$
\prod_{\text {forward }} \hat{\eta}_{\text {forward }}=\prod_{\text {back }} \hat{\eta}_{\text {back }}
$$

If the velocity distribution is Maxwellian, this gives:

$$
\hat{\eta}_{x}=\frac{n_{x}}{g_{x}} \frac{h^{3}}{\left(2 \pi m_{x} k T\right)^{3 / 2}} \mathrm{e}^{-E_{\text {kin }} / k T}
$$

where $g$ is the statistical weight of the state and $n / g:=\eta$. For electrons holds $g=2$, for excited states usually holds $g=2 j+1=2 n^{2}$.

With this one finds for the Boltzmann balance, $\mathrm{X}_{p}+\mathrm{e}^{-} \rightleftarrows \mathrm{X}_{1}+\mathrm{e}^{-}+\left(E_{1 p}\right)$ :

$$
\frac{n_{p}^{\mathrm{B}}}{n_{1}}=\frac{g_{p}}{g_{1}} \exp \left(\frac{E_{p}-E_{1}}{k T_{\mathrm{e}}}\right)
$$

And for the Saha balance, $\mathrm{X}_{p}+\mathrm{e}^{-}+\left(E_{p i}\right) \rightleftarrows \mathrm{X}_{1}^{+}+2 \mathrm{e}^{-}$:

$$
\frac{n_{p}^{\mathrm{S}}}{g_{p}}=\frac{n_{1}^{+}}{g_{1}^{+}} \frac{n_{\mathrm{e}}}{g_{\mathrm{e}}} \frac{h^{3}}{\left(2 \pi m_{\mathrm{e}} k T_{\mathrm{e}}\right)^{3 / 2}} \exp \left(\frac{E_{p i}}{k T_{\mathrm{e}}}\right)
$$

Because the number of particles on the left-hand side and right-hand side of the equation is different, a factor $g / V_{\mathrm{e}}$ remains. This factor causes the Saha-jump.
From microscopic reversibility one can derive that for the rate coefficients $K(p, q, T):=\langle\sigma v\rangle_{p q}$ holds:

$$
K(q, p, T)=\frac{g_{p}}{g_{q}} K(p, q, T) \exp \left(\frac{\Delta E_{p q}}{k T}\right)
$$

### 11.5 Inelastic collisions

### 11.5.1 Types of collisions

The kinetic energy can be split in a part of and a part in the centre of mass system. The energy in the centre of mass system is available for reactions. This energy is given by

$$
E=\frac{m_{1} m_{2}\left(v_{1}-v_{2}\right)^{2}}{2\left(m_{1}+m_{2}\right)}
$$

Some types of inelastic collisions important for plasma physics are:

1. Excitation: $\mathrm{A}_{p}+\mathrm{e}^{-} \rightleftarrows \mathrm{A}_{q}+\mathrm{e}^{-}$
2. Decay: $\mathrm{A}_{q} \rightleftarrows \mathrm{~A}_{p}+h f$
3. Ionisation and 3-particles recombination: $\mathrm{A}_{p}+\mathrm{e}^{-} \rightleftarrows \mathrm{A}^{+}+2 \mathrm{e}^{-}$
4. radiative recombination: $\mathrm{A}^{+}+\mathrm{e}^{-} \rightleftarrows \mathrm{A}_{p}+h f$
5. Stimulated emission: $\mathrm{A}_{q}+h f \rightarrow \mathrm{~A}_{p}+2 h f$
6. Associative ionisation: $\mathrm{A}^{* *}+\mathrm{B} \rightleftarrows \mathrm{AB}^{+}+\mathrm{e}^{-}$
7. Penning ionisation: b.v. $\mathrm{Ne}^{*}+\mathrm{Ar} \rightleftarrows \mathrm{Ar}^{+}+\mathrm{Ne}+\mathrm{e}^{-}$
8. Charge transfer: $\mathrm{A}^{+}+\mathrm{B} \rightleftarrows \mathrm{A}+\mathrm{B}^{+}$
9. Resonant charge transfer: $\mathrm{A}^{+}+\mathrm{A} \rightleftarrows \mathrm{A}+\mathrm{A}^{+}$

### 11.5.2 Cross sections

Collisions between an electron and an atom can be approximated by a collision between an electron and one of the electrons of that atom. This results in

$$
\frac{d \sigma}{d(\Delta E)}=\frac{\pi Z^{2} e^{4}}{\left(4 \pi \varepsilon_{0}\right)^{2} E(\Delta E)^{2}}
$$

Then follows for the transition $p \rightarrow q: \sigma_{p q}(E)=\frac{\pi Z^{2} e^{4} \Delta E_{q, q+1}}{\left(4 \pi \varepsilon_{0}\right)^{2} E(\Delta E)_{p q}^{2}}$
For ionization from state $p$ holds to a good approximation: $\sigma_{p}=4 \pi a_{0}^{2} R y\left(\frac{1}{E_{p}}-\frac{1}{E}\right) \ln \left(\frac{1.25 \beta E}{E_{p}}\right)$
For resonant charge transfer holds: $\sigma_{\text {ex }}=\frac{A[1-B \ln (E)]^{2}}{1+C E^{3.3}}$

### 11.6 Radiation

In equilibrium holds for radiation processes:

$$
\underbrace{n_{p} A_{p q}}_{\text {emission }}+\underbrace{n_{p} B_{p q} \rho(\nu, T)}_{\text {stimulated emission }}=\underbrace{n_{q} B_{q p} \rho(\nu, T)}_{\text {absorption }}
$$

Here, $A_{p q}$ is the matrix element of the transition $p \rightarrow q$, and is given by:

$$
A_{p q}=\frac{8 \pi^{2} e^{2} \nu^{3}\left|r_{p q}\right|^{2}}{3 \hbar \varepsilon_{0} c^{3}} \text { with } r_{p q}=\left\langle\psi_{p}\right| \vec{r}\left|\psi_{q}\right\rangle
$$

For hydrogenic atoms holds: $A_{p}=1.58 \cdot 10^{8} Z^{4} p^{-4.5}$, with $A_{p}=1 / \tau_{p}=\sum_{q} A_{p q}$. The intensity $I$ of a line is given by $I_{p q}=h f A_{p q} n_{p} / 4 \pi$. The Einstein coefficients $B$ are given by:

$$
B_{p q}=\frac{c^{3} A_{p q}}{8 \pi h \nu^{3}} \text { and } \frac{B_{p q}}{B_{q p}}=\frac{g_{q}}{g_{p}}
$$

A spectral line is broadened by several mechanisms:

1. Because the states have a finite life time. The natural life time of a state $p$ is given by $\tau_{p}=1 / \sum_{q} A_{p q}$.

From the uncertainty relation then follows: $\Delta(h \nu) \cdot \tau_{p}=\frac{1}{2} \hbar$, this gives

$$
\Delta \nu=\frac{1}{4 \pi \tau_{p}}=\frac{\sum_{q} A_{p q}}{4 \pi}
$$

The natural line width is usually $\ll$ than the broadening due to the following two mechanisms:
2. The Doppler broadening is caused by the thermal motion of the particles:

$$
\frac{\Delta \lambda}{\lambda}=\frac{2}{c} \sqrt{\frac{2 \ln (2) k T_{\mathrm{i}}}{m_{\mathrm{i}}}}
$$

This broadening results in a Gaussian line profile:
$k_{\nu}=k_{0} \exp \left(-\left[2 \sqrt{\ln 2}\left(\nu-\nu_{0}\right) / \Delta \nu_{D}\right]^{2}\right)$, with $k$ the coefficient of absorption or emission.
3. The Stark broadening is caused by the electric field of the electrons:

$$
\Delta \lambda_{1 / 2}=\left[\frac{n_{\mathrm{e}}}{C\left(n_{\mathrm{e}}, T_{\mathrm{e}}\right)}\right]^{2 / 3}
$$

with for the $\mathrm{H}-\beta$ line: $C\left(n_{\mathrm{e}}, T_{\mathrm{e}}\right) \approx 3 \cdot 10^{14} \AA^{-3 / 2} \mathrm{~cm}^{-3}$.
The natural broadening and the Stark broadening result in a Lorentz profile of a spectral line:
$k_{\nu}=\frac{1}{2} k_{0} \Delta \nu_{L} /\left[\left(\frac{1}{2} \Delta \nu_{L}\right)^{2}+\left(\nu-\nu_{0}\right)^{2}\right]$. The total line shape is a convolution of the Gauss- and Lorentz profile and is called a Voigt profile.
The number of transitions $p \rightarrow q$ is given by $n_{p} B_{p q} \rho$ and by $n_{p} n_{h f}\left\langle\sigma_{\mathrm{a}} c\right\rangle=n_{p}(\rho d \nu / h \nu) \sigma_{\mathrm{a}} c$ where $d \nu$ is the line width. Then follows for the cross section of absorption processes: $\sigma_{\mathrm{a}}=B_{p q} h \nu / c d \nu$.

The background radiation in a plasma originates from two processes:

1. Free-Bound radiation, originating from radiative recombination. The emission is given by:

$$
\varepsilon_{f b}=\frac{C_{1}}{\lambda^{2}} \frac{z_{\mathrm{i}} n_{\mathrm{i}} n_{\mathrm{e}}}{\sqrt{k T_{\mathrm{e}}}}\left[1-\exp \left(-\frac{h c}{\lambda k T_{\mathrm{e}}}\right)\right] \xi_{f b}\left(\lambda, T_{\mathrm{e}}\right)
$$

with $C_{1}=1.63 \cdot 10^{-43} \mathrm{Wm}^{4} \mathrm{~K}^{1 / 2} \mathrm{sr}^{-1}$ and $\xi$ the Biberman factor.
2. Free-free radiation, originating from the acceleration of particles in the EM-field of other particles:

$$
\varepsilon_{f f}=\frac{C_{1}}{\lambda^{2}} \frac{z_{\mathrm{i}} n_{\mathrm{i}} n_{\mathrm{e}}}{\sqrt{k T_{\mathrm{e}}}} \exp \left(-\frac{h c}{\lambda k T_{\mathrm{e}}}\right) \xi_{f f}\left(\lambda, T_{\mathrm{e}}\right)
$$

### 11.7 The Boltzmann transport equation

It is assumed that there exists a distribution function $F$ for the plasma so that

$$
F(\vec{r}, \vec{v}, t)=F_{r}(\vec{r}, t) \cdot F_{v}(\vec{v}, t)=F_{1}(x, t) F_{2}(y, t) F_{3}(z, t) F_{4}\left(v_{x}, t\right) F_{5}\left(v_{y}, t\right) F_{6}\left(v_{z}, t\right)
$$

Then the BTE is: $\frac{d F}{d t}=\frac{\partial F}{\partial t}+\nabla_{r} \cdot(F \vec{v})+\nabla_{v} \cdot(F \vec{a})=\left(\frac{\partial F}{\partial t}\right)_{\text {coll-rad }}$
Assuming that $v$ does not depend on $r$ and $a_{i}$ does not depend on $v_{i}$, holds $\nabla_{r} \cdot(F \vec{v})=\vec{v} \cdot \nabla F$ and $\nabla_{v} \cdot(F \vec{a})=$ $\vec{a} \cdot \nabla_{v} F$. This is also true in magnetic fields because $\partial a_{i} / \partial x_{i}=0$. The velocity is separated in a thermal velocity $\vec{v}_{\mathrm{t}}$ and a drift velocity $\vec{w}$. The total density is given by $n=\int F d \vec{v}$ and $\int \vec{v} F d \vec{v}=n \vec{w}$.

The balance equations can be derived by means of the moment method:

1. Mass balance: $\int(\mathrm{BTE}) d \vec{v} \Rightarrow \frac{\partial n}{\partial t}+\nabla \cdot(n \vec{w})=\left(\frac{\partial n}{\partial t}\right)_{\mathrm{cr}}$
2. Momentum balance: $\int(\mathrm{BTE}) m \vec{v} d \vec{v} \Rightarrow m n \frac{d \vec{w}}{d t}+\nabla \mathbf{T}^{\prime}+\nabla p=m n\langle\vec{a}\rangle+\vec{R}$
3. Energy balance: $\int(\mathrm{BTE}) m v^{2} d \vec{v} \Rightarrow \frac{3}{2} \frac{d p}{d t}+\frac{5}{2} p \nabla \cdot \vec{w}+\nabla \cdot \vec{q}=Q$

Here, $\langle\vec{a}\rangle=e / m(\vec{E}+\vec{w} \times \vec{B})$ is the average acceleration, $\vec{q}=\frac{1}{2} n m\left\langle\vec{v}_{\mathrm{t}}^{2} \vec{v}_{\mathrm{t}}\right\rangle$ the heat flow, $Q=\int_{\mathrm{c}} \frac{m v_{\mathrm{t}}^{2}}{r}\left(\frac{\partial F}{\partial t}\right)_{\mathrm{cr}} d \vec{v}$ the source term for energy production, $\vec{R}$ is a friction term and $p=n k T$ the pressure.
A thermodynamic derivation gives for the total pressure: $p=n k T=\sum_{i} p_{i}-\frac{e^{2}\left(n_{\mathrm{e}}+z_{\mathrm{i}} n_{\mathrm{i}}\right)}{24 \pi \varepsilon_{0} \lambda_{\mathrm{D}}}$
For the electrical conductance in a plasma follows from the momentum balance, if $w_{\mathrm{e}} \gg w_{\mathrm{i}}$ :

$$
\eta \vec{J}=\vec{E}-\frac{\vec{J} \times \vec{B}+\nabla p_{\mathrm{e}}}{e n_{\mathrm{e}}}
$$

In a plasma where only elastic e-a collisions are important the equilibrium energy distribution function is the Druyvesteyn distribution:

$$
N(E) d E=C n_{\mathrm{e}}\left(\frac{E}{E_{0}}\right)^{3 / 2} \exp \left[-\frac{3 m_{\mathrm{e}}}{m_{0}}\left(\frac{E}{E_{0}}\right)^{2}\right] d E
$$

with $E_{0}=e E \lambda_{\mathrm{v}}=e E / n \sigma$.

### 11.8 Collision-radiative models

These models are first-moment equations for excited states. One assumes the Quasi-steady-state solution is valid, where $\forall_{p>1}\left[\left(\partial n_{p} / \partial t=0\right) \wedge\left(\nabla \cdot\left(n_{p} \vec{w}_{p}\right)=0\right)\right]$. This results in:

$$
\left(\frac{\partial n_{p>1}}{\partial t}\right)_{\mathrm{cr}}=0, \frac{\partial n_{1}}{\partial t}+\nabla \cdot\left(n_{1} \vec{w}_{1}\right)=\left(\frac{\partial n_{1}}{\partial t}\right)_{\mathrm{cr}}, \frac{\partial n_{\mathrm{i}}}{\partial t}+\nabla \cdot\left(n_{\mathrm{i}} \vec{w}_{\mathrm{i}}\right)=\left(\frac{\partial n_{\mathrm{i}}}{\partial t}\right)_{\mathrm{cr}}
$$

with solutions $n_{p}=r_{p}^{0} n_{p}^{\mathrm{S}}+r_{p}^{1} n_{p}^{\mathrm{B}}=b_{p} n_{p}^{\mathrm{S}}$. Further holds for all collision-dominated levels that $\delta b_{p}:=b_{p}-1=$ $b_{0} p_{\text {eff }}^{-x}$ with $p_{\text {eff }}=\sqrt{R y / E_{p \mathrm{i}}}$ and $5 \leq x \leq 6$. For systems in ESP, where only collisional (de)excitation between levels $p$ and $p \pm 1$ is taken into account holds $x=6$. Even in plasma's far from equilibrium the excited levels will eventually reach ESP, so from a certain level up the level densities can be calculated.

To find the population densities of the lower levels in the stationary case one has to start with a macroscopic equilibrium:

Number of populating processes of level $p=$ Number of depopulating processes of level $p$,
When this is expanded it becomes:

$$
\begin{aligned}
& \underbrace{n_{\mathrm{e}} \sum_{q<p} n_{q} K_{q p}}_{\text {coll. excit. }}+\underbrace{n_{\mathrm{e}} \sum_{q>p} n_{q} K_{q p}}_{\text {coll. deexcit. }}+\underbrace{\sum_{q>p} n_{q} A_{q p}}_{\text {rad. deex. to }}+\underbrace{n_{e}^{2} n_{\mathrm{i}} K_{+p}}_{\text {coll. recomb. }}+\underbrace{n_{\mathrm{e}} n_{\mathrm{i}} \alpha_{\mathrm{rad}}}_{\text {rad. recomb }}= \\
& \quad \underbrace{n_{\mathrm{e}} n_{p} \sum_{q<p} K_{p q}}_{\text {coll. deexcit. }}+\underbrace{n_{\mathrm{e}} n_{p} \sum_{q>p} K_{p q}}_{\text {coll. excit. }}+\underbrace{n_{p} \sum_{q<p} A_{p q}}_{\text {rad. deex. from }}+\underbrace{n_{\mathrm{e}} n_{p} K_{p+}}_{\text {coll. ion. }}
\end{aligned}
$$

### 11.9 Waves in plasma's

Interaction of electromagnetic waves in plasma's results in scattering and absorption of energy. For electromagnetic waves with complex wave number $k=\omega(n+i \kappa) / c$ in one dimension one finds: $E_{x}=E_{0} \mathrm{e}^{-\kappa \omega x / c} \cos [\omega(t-n x / c)]$. The refractive index $n$ is given by:

$$
n=c \frac{k}{\omega}=\frac{c}{v_{\mathrm{f}}}=\sqrt{1-\frac{\omega_{\mathrm{p}}^{2}}{\omega^{2}}}
$$

For disturbances in the $z$-direction in a cold, homogeneous, magnetized plasma: $\vec{B}=B_{0} \vec{e}_{z}+\overrightarrow{\hat{B}} \mathrm{e}^{i(k z-\omega t)}$ and $n=n_{0}+\hat{n} \mathrm{e}^{i(k z-\omega t)}$ (external $E$ fields are screened) follows, with the definitions $\alpha=\omega_{\mathrm{p}} / \omega$ and $\beta=\Omega / \omega$ and $\omega_{\mathrm{p}}^{2}=\omega_{\mathrm{pi}}^{2}+\omega_{\mathrm{pe}}^{2}$ :

$$
\vec{J}=\overrightarrow{\vec{\sigma}} \vec{E} \text {, with } \overrightarrow{\vec{\sigma}}=i \varepsilon_{0} \omega \sum_{s} \alpha_{s}^{2}\left(\begin{array}{ccc}
\frac{1}{1-\beta_{s}^{2}} & \frac{-i \beta_{s}}{1-\beta_{s}^{2}} & 0 \\
\frac{i \beta_{s}}{1-\beta_{s}^{2}} & \frac{1}{1-\beta_{s}^{2}} & 0 \\
0 & 0 & 1
\end{array}\right)
$$

where the sum is taken over particle species $s$. The dielectric tensor $\mathcal{E}$, with property:

$$
\vec{k} \cdot(\overrightarrow{\overrightarrow{\mathcal{E}}} \cdot \vec{E})=0
$$

is given by $\overrightarrow{\mathcal{E}}=\overrightarrow{\vec{I}}-\vec{\sigma} / i \varepsilon_{0} \omega$.
With the definitions $S=1-\sum_{s} \frac{\alpha_{s}^{2}}{1-\beta_{s}^{2}}, \quad D=\sum_{s} \frac{\alpha_{s}^{2} \beta_{s}}{1-\beta_{s}^{2}}, \quad P=1-\sum_{s} \alpha_{s}^{2}$
follows:

$$
\overrightarrow{\overrightarrow{\mathcal{E}}}=\left(\begin{array}{ccc}
S & -i D & 0 \\
i D & S & 0 \\
0 & 0 & P
\end{array}\right)
$$

The eigenvalues of this hermitian matrix are $R=S+D, L=S-D, \lambda_{3}=P$, with eigenvectors $\vec{e}_{\mathrm{r}}=$ $\frac{1}{2} \sqrt{2}(1, i, 0), \vec{e}_{1}=\frac{1}{2} \sqrt{2}(1,-i, 0)$ and $\vec{e}_{3}=(0,0,1) . \vec{e}_{\mathrm{r}}$ is connected with a right rotating field for which $i E_{x} / E_{y}=1$ and $\vec{e}_{1}$ is connected with a left rotating field for which $i E_{x} / E_{y}=-1$. When $k$ makes an angle $\theta$ with $\vec{B}$ one finds:

$$
\tan ^{2}(\theta)=\frac{P\left(n^{2}-R\right)\left(n^{2}-L\right)}{S\left(n^{2}-R L / S\right)\left(n^{2}-P\right)}
$$

where $n$ is the refractive index. From this the following solutions can be obtained:

## A. $\theta=0$ : transmission in the $\boldsymbol{z}$-direction.

1. $P=0: E_{x}=E_{y}=0$. This describes a longitudinal linear polarized wave.
2. $n^{2}=L$ : a left, circular polarized wave.
3. $n^{2}=R$ : a right, circular polarized wave.

## B. $\theta=\pi / 2$ : transmission $\perp$ the $B$-field.

1. $n^{2}=P$ : the ordinary mode: $E_{x}=E_{y}=0$. This is a transversal linear polarized wave.
2. $n^{2}=R L / S$ : the extraordinary mode: $i E_{x} / E_{y}=-D / S$, an elliptical polarized wave.

Resonance frequencies are frequencies for which $n^{2} \rightarrow \infty$, so $v_{\mathrm{f}}=0$. For these holds: $\tan (\theta)=-P / S$. For $R \rightarrow \infty$ this gives the electron cyclotron resonance frequency $\omega=\Omega_{\mathrm{e}}$, for $L \rightarrow \infty$ the ion cyclotron resonance frequency $\omega=\Omega_{\mathrm{i}}$ and for $S=0$ holds for the extraordinary mode:

$$
\alpha^{2}\left(1-\frac{m_{\mathrm{i}}}{m_{\mathrm{e}}} \frac{\Omega_{\mathrm{i}}^{2}}{\omega^{2}}\right)=\left(1-\frac{m_{\mathrm{i}}^{2}}{m_{\mathrm{e}}^{2}} \frac{\Omega_{\mathrm{i}}^{2}}{\omega^{2}}\right)\left(1-\frac{\Omega_{\mathrm{i}}^{2}}{\omega^{2}}\right)
$$

Cut-off frequencies are frequencies for which $n^{2}=0$, so $v_{\mathrm{f}} \rightarrow \infty$. For these holds: $P=0$ or $R=0$ or $L=0$. In the case that $\beta^{2} \gg 1$ one finds Alfvén waves propagating parallel to the field lines. With the Alfvén velocity

$$
v_{\mathrm{A}}=\frac{\Omega_{\mathrm{e}} \Omega_{\mathrm{i}}}{\omega_{\mathrm{pe}}^{2}+\omega_{\mathrm{pi}}^{2}} c^{2}
$$

follows: $n=\sqrt{1+c / v_{\mathrm{A}}}$, and in case $v_{\mathrm{A}} \ll c: \omega=k v_{\mathrm{A}}$.

## Chapter 12

## Solid state physics

### 12.1 Crystal structure

A lattice is defined by the 3 translation vectors $\vec{a}_{i}$, so that the atomic composition looks the same from each point $\vec{r}$ and $\vec{r}^{\prime}=\vec{r}+\vec{T}$, where $\vec{T}$ is a translation vector given by: $\vec{T}=u_{1} \vec{a}_{1}+u_{2} \vec{a}_{2}+u_{3} \vec{a}_{3}$ with $u_{i} \in I N$. A lattice can be constructed from primitive cells. As a primitive cell one can take a parallellepiped, with volume

$$
V_{\mathrm{cell}}=\left|\vec{a}_{1} \cdot\left(\vec{a}_{2} \times \vec{a}_{3}\right)\right|
$$

Because a lattice has a periodical structure the physical properties which are connected with the lattice have the same periodicity (neglecting boundary effects):

$$
n_{\mathrm{e}}(\vec{r}+\vec{T})=n_{\mathrm{e}}(\vec{r})
$$

This periodicity is suitable to use Fourier analysis: $n(\vec{r})$ is expanded as:

$$
n(\vec{r})=\sum_{G} n_{G} \exp (i \vec{G} \cdot \vec{r})
$$

with

$$
n_{G}=\frac{1}{V_{\text {cell }}} \iiint_{\text {cell }} n(\vec{r}) \exp (-i \vec{G} \cdot \vec{r}) d V
$$

$\vec{G}$ is the reciprocal lattice vector. If $\vec{G}$ is written as $\vec{G}=v_{1} \vec{b}_{1}+v_{2} \vec{b}_{2}+v_{3} \vec{b}_{3}$ with $v_{i} \in I N$, it follows for the vectors $\vec{b}_{i}$, cyclically:

$$
\vec{b}_{i}=2 \pi \frac{\vec{a}_{i+1} \times \vec{a}_{i+2}}{\vec{a}_{i} \cdot\left(\vec{a}_{i+1} \times \vec{a}_{i+2}\right)}
$$

The set of $\vec{G}$-vectors determines the Röntgen diffractions: a maximum in the reflected radiation occurs if: $\Delta \vec{k}=\vec{G}$ with $\Delta \vec{k}=\vec{k}-\vec{k}^{\prime}$. So: $2 \vec{k} \cdot \vec{G}=G^{2}$. From this follows for parallel lattice planes (Bragg reflection) that for the maxima holds: $2 d \sin (\theta)=n \lambda$.
The Brillouin zone is defined as a Wigner-Seitz cell in the reciprocal lattice.

### 12.2 Crystal binding

A distinction can be made between 4 binding types:

1. Van der Waals bond
2. Ion bond
3. Covalent or homopolar bond
4. Metalic bond.

For the ion binding of NaCl the energy per molecule is calculated by:

$$
\mathrm{E}=\text { cohesive energy }(\mathrm{NaCl})-\text { ionization energy }(\mathrm{Na})+\text { electron affinity }(\mathrm{Cl})
$$

The interaction in a covalent bond depends on the relative spin orientations of the electrons constituing the bond. The potential energy for two parallel spins is higher than the potential energy for two antiparallel spins. Furthermore the potential energy for two parallel spins has sometimes no minimum. In that case binding is not possible.

### 12.3 Crystal vibrations

### 12.3.1 A lattice with one type of atoms

In this model for crystal vibrations only nearest-neighbour interactions are taken into account. The force on atom $s$ with mass $M$ can then be written as:

$$
F_{s}=M \frac{d^{2} u_{s}}{d t^{2}}=C\left(u_{s+1}-u_{s}\right)+C\left(u_{s-1}-u_{s}\right)
$$

Assuming that all solutions have the same time-dependence $\exp (-i \omega t)$ this results in:

$$
-M \omega^{2} u_{s}=C\left(u_{s+1}+u_{s-1}-2 u_{s}\right)
$$

Further it is postulated that: $u_{s \pm 1}=u \exp (i s K a) \exp ( \pm i K a)$.
This gives: $u_{s}=\exp (i K s a)$. Substituting the later two equations in the fist results in a system of linear equations, which has only a solution if their determinant is 0 . This gives:

$$
\omega^{2}=\frac{4 C}{M} \sin ^{2}\left(\frac{1}{2} K a\right)
$$

Only vibrations with a wavelength within the first Brillouin Zone have a physical significance. This requires that $-\pi<K a \leq \pi$.

The group velocity of these vibrations is given by:

$$
v_{\mathrm{g}}=\frac{d \omega}{d K}=\sqrt{\frac{C a^{2}}{M}} \cos \left(\frac{1}{2} K a\right)
$$

and is 0 on the edge of a Brillouin Zone. Here, there is a standing wave.

### 12.3.2 A lattice with two types of atoms

Now the solutions are:

$$
\omega^{2}=C\left(\frac{1}{M_{1}}+\frac{1}{M_{2}}\right) \pm C \sqrt{\left(\frac{1}{M_{1}}+\frac{1}{M_{2}}\right)^{2}-\frac{4 \sin ^{2}(K a)}{M_{1} M_{2}}}
$$

Connected with each value of $K$ are two values of $\omega$, as can be seen in the graph. The upper line describes the optical branch, the lower line the acoustical branch. In the optical branch, both types of ions oscillate in opposite phases, in the acoustical branch they oscillate in the same phase. This results in a much
 larger induced dipole moment for optical oscillations, and also a stronger emission and absorption of radiation. Furthermore each branch has 3 polarization directions, one longitudinal and two transversal.

### 12.3.3 Phonons

The quantum mechanical excitation of a crystal vibration with an energy $\hbar \omega$ is called a phonon. Phonons can be viewed as quasi-particles: with collisions, they behave as particles with momentum $\hbar K$. Their total momentum is 0 . When they collide, their momentum need not be conserved: for a normal process holds: $K_{1}+K_{2}=K_{3}$, for an umklapp process holds: $K_{1}+K_{2}=K_{3}+G$. Because phonons have no spin they behave like bosons.

### 12.3.4 Thermal heat capacity

The total energy of the crystal vibrations can be calculated by multiplying each mode with its energy and sum over all branches $K$ and polarizations $P$ :

$$
U=\sum_{K} \sum_{P} \hbar \omega\left\langle n_{k, p}\right\rangle=\sum_{\lambda} \int D_{\lambda}(\omega) \frac{\hbar \omega}{\exp (\hbar \omega / k T)-1} d \omega
$$

for a given polarization $\lambda$. The thermal heat capacity is then:

$$
C_{\text {lattice }}=\frac{\partial U}{\partial T}=k \sum_{\lambda} \int D(\omega) \frac{(\hbar \omega / k T)^{2} \exp (\hbar \omega / k T)}{(\exp (\hbar \omega / k T)-1)^{2}} d \omega
$$

The dispersion relation in one dimension is given by:

$$
D(\omega) d \omega=\frac{L}{\pi} \frac{d K}{d \omega} d \omega=\frac{L}{\pi} \frac{d \omega}{v_{\mathrm{g}}}
$$

In three dimensions one applies periodic boundary conditions to a cube with $N^{3}$ primitive cells and a volume $L^{3}: \exp \left(i\left(K_{x} x+K_{y} y+K_{z} z\right)\right) \equiv \exp \left(i\left(K_{x}(x+L)+K_{y}(y+L)+K_{z}(z+L)\right)\right)$.
Because $\exp (2 \pi i)=1$ this is only possible if:

$$
K_{x}, K_{y}, K_{z}=0 ; \pm \frac{2 \pi}{L} ; \pm \frac{4 \pi}{L} ; \pm \frac{6 \pi}{L} ; \ldots \pm \frac{2 N \pi}{L}
$$

So there is only one allowed value of $\vec{K}$ per volume $(2 \pi / L)^{3}$ in $K$-space, or:

$$
\left(\frac{L}{2 \pi}\right)^{3}=\frac{V}{8 \pi^{3}}
$$

allowed $\vec{K}$-values per unit volume in $\vec{K}$-space, for each polarization and each branch. The total number of states with a wave vector $<K$ is:

$$
N=\left(\frac{L}{2 \pi}\right)^{3} \frac{4 \pi K^{3}}{3}
$$

for each polarization. The density of states for each polarization is, according to the Einstein model:

$$
D(\omega)=\frac{d N}{d \omega}=\left(\frac{V K^{2}}{2 \pi^{2}}\right) \frac{d K}{d \omega}=\frac{V}{8 \pi^{3}} \iint \frac{d A_{\omega}}{v_{\mathrm{g}}}
$$

The Debye model for thermal heat capacities is a low-temperature approximation which is valid up to $\approx 50 \mathrm{~K}$. Here, only the acoustic phonons are taken into account (3 polarizations), and one assumes that $v=\omega K$, independent of the polarization. From this follows: $D(\omega)=V \omega^{2} / 2 \pi^{2} v^{3}$, where $v$ is the speed of sound. This gives:

$$
U=3 \int D(\omega)\langle n\rangle \hbar \omega d \omega=\int_{0}^{\omega_{D}} \frac{V \omega^{2}}{2 \pi^{2} v^{3}} \frac{\hbar \omega}{\exp (\hbar \omega / k T)-1} d \omega=\frac{3 V k^{2} T^{4}}{2 \pi^{2} v^{3} \hbar^{3}} \int_{0}^{x_{\mathrm{D}}} \frac{x^{3} d x}{\mathrm{e}^{x}-1}
$$

Here, $x_{D}=\hbar \omega_{\mathrm{D}} / k T=\theta_{\mathrm{D}} / T . \theta_{\mathrm{D}}$ is the Debye temperature and is defined by:

$$
\theta_{\mathrm{D}}=\frac{\hbar v}{k}\left(\frac{6 \pi^{2} N}{V}\right)^{1 / 3}
$$

where $N$ is the number of primitive cells. Because $x_{\mathrm{D}} \rightarrow \infty$ for $T \rightarrow 0$ it follows from this:

$$
U=9 N k T\left(\frac{T}{\theta_{\mathrm{D}}}\right)^{3} \int_{0}^{\infty} \frac{x^{3} d x}{\mathrm{e}^{x}-1}=\frac{3 \pi^{4} N k T^{4}}{5 \theta_{\mathrm{D}}} \sim T^{4} \quad \text { and } \quad C_{V}=\frac{12 \pi^{4} N k T^{3}}{5 \theta_{\mathrm{D}}^{3}} \sim T^{3}
$$

In the Einstein model for the thermal heat capacity one considers only phonons at one frequency, an approximation for optical phonons.

### 12.4 Magnetic field in the solid state

The following graph shows the magnetization versus fieldstrength for different types of magnetism:

$$
M_{\text {sat }}^{M}
$$

### 12.4.1 Dielectrics

The quantum mechanical origin of diamagnetism is the Larmorprecession of the spin of the electron. Starting with a circular electron orbit in an atom with two electrons, there is a Coulomb force $F_{\mathrm{c}}$ and a magnetic force on each electron. If the magnetic part of the force is not strong enough to significantly deform the orbit holds:

$$
\omega^{2}=\frac{F_{\mathrm{c}}(r)}{m r} \pm \frac{e B}{m} \omega=\omega_{0}^{2} \pm \frac{e B}{m}\left(\omega_{0}+\delta\right) \Rightarrow \omega=\sqrt{\left(\omega_{0} \pm \frac{e B}{2 m}\right)^{2}+\cdots} \approx \omega_{0} \pm \frac{e B}{2 m}=\omega_{0} \pm \omega_{\mathrm{L}}
$$

Here, $\omega_{\mathrm{L}}$ is the Larmor frequency. One electron is accelerated, the other decelerated. Hence there is a net circular current which results in a magnetic moment $\vec{\mu}$. The circular current is given by $I=-Z e \omega_{\mathrm{L}} / 2 \pi$, and $\langle\mu\rangle=I A=I \pi\left\langle\rho^{2}\right\rangle=\frac{2}{3} I \pi\left\langle r^{2}\right\rangle$. If $N$ is the number of atoms in the crystal it follows for the susceptibility, with $\vec{M}=\vec{\mu} N$ :

$$
\chi=\frac{\mu_{0} M}{B}=-\frac{\mu_{0} N Z e^{2}}{6 m}\left\langle r^{2}\right\rangle
$$

### 12.4.2 Paramagnetism

Starting with the splitting of energy levels in a weak magnetic field: $\Delta U_{m}-\vec{\mu} \cdot \vec{B}=m_{J} g \mu_{\mathrm{B}} B$, and with a distribution $f_{m} \sim \exp \left(-\Delta U_{m} / k T\right)$, one finds for the average magnetic moment $\langle\mu\rangle=\sum f_{m} \mu / \sum f_{m}$. After linearization and because $\sum m_{J}=0, \sum J=2 J+1$ and $\sum m_{J}^{2}=\frac{2}{3} J(J+1)\left(J+\frac{1}{2}\right)$ it follows that:

$$
\chi_{p}=\frac{\mu_{0} M}{B}=\frac{\mu_{0} N\langle\mu\rangle}{B}=\frac{\mu_{0} J(J+1) g^{2} \mu_{\mathrm{B}}^{2} N}{3 k T}
$$

This is the Curie law, $\chi_{p} \sim 1 / T$.

### 12.4.3 Ferromagnetism

A ferromagnet behaves like a paramagnet above a critical temperature $T_{\mathrm{c}}$. To describe ferromagnetism a field $B_{E}$ parallel with $M$ is postulated: $\vec{B}_{E}=\lambda \mu_{0} \vec{M}$. From there the treatment is analogous to the paramagnetic case:

$$
\mu_{0} M=\chi_{p}\left(B_{\mathrm{a}}+B_{E}\right)=\chi_{p}\left(B_{\mathrm{a}}+\lambda \mu_{0} M\right)=\mu_{0}\left(1-\lambda \frac{C}{T}\right) M
$$

From this follows for a ferromagnet: $\chi_{F}=\frac{\mu_{0} M}{B_{\mathrm{a}}}=\frac{C}{T-T_{\mathrm{c}}}$ which is Weiss-Curie's law.
If $B_{E}$ is estimated this way it results in values of about 1000 T . This is clearly unrealistic and suggests another mechanism. A quantum mechanical approach from Heisenberg postulates an interaction between two neighbouring atoms: $U=-2 J \vec{S}_{i} \cdot \vec{S}_{j} \equiv-\vec{\mu} \cdot \vec{B}_{E} . J$ is an overlap integral given by: $J=3 k T_{\mathrm{c}} / 2 z S(S+1)$, with $z$ the number of neighbours. A distinction between 2 cases can now be made:

1. $J>0: S_{i}$ and $S_{j}$ become parallel: the material is a ferromagnet.
2. $J<0: S_{i}$ and $S_{j}$ become antiparallel: the material is an antiferromagnet.

Heisenberg's theory predicts quantized spin waves: magnons. Starting from a model with only nearest neighbouring atoms interacting one can write:

$$
U=-2 J \vec{S}_{p} \cdot\left(\vec{S}_{p-1}+\vec{S}_{p+1}\right) \approx \vec{\mu}_{p} \cdot \vec{B}_{p} \quad \text { with } \quad \vec{B}_{p}=\frac{-2 J}{g \mu_{\mathrm{B}}}\left(\vec{S}_{p-1}+\vec{S}_{p+1}\right)
$$

The equation of motion for the magnons becomes: $\frac{d \vec{S}}{d t}=\frac{2 J}{\hbar} \vec{S}_{p} \times\left(\vec{S}_{p-1}+\vec{S}_{p+1}\right)$
From here the treatment is analogous to phonons: postulate traveling waves of the type $\vec{S}_{p}=\vec{u} \exp (i) p k a-$ $\omega t)$ ). This results in a system of linear equations with solution:

$$
\hbar \omega=4 J S(1-\cos (k a))
$$

### 12.5 Free electron Fermi gas

### 12.5.1 Thermal heat capacity

The solution with period $L$ of the one-dimensional Schrödinger equation is: $\psi_{n}(x)=A \sin (2 \pi x / \lambda n)$ with $n \lambda_{n}=2 L$. From this follows

$$
E=\frac{\hbar^{2}}{2 m}\left(\frac{n \pi}{L}\right)^{2}
$$

In a linear lattice the only important quantum numbers are $n$ and $m_{s}$. The Fermi level is the uppermost filled level in the ground state, which has the Fermi-energy $E_{\mathrm{F}}$. If $n_{\mathrm{F}}$ is the quantum number of the Fermi level, it can be expressed as: $2 n_{\mathrm{F}}=N$ so $E_{\mathrm{F}}=\hbar^{2} \pi^{2} N^{2} / 8 m L$. In 3 dimensions holds:

$$
k_{\mathrm{F}}=\left(\frac{3 \pi^{2} N}{V}\right)^{1 / 3} \text { and } E_{\mathrm{F}}=\frac{\hbar^{2}}{2 m}\left(\frac{3 \pi^{2} N}{V}\right)^{2 / 3}
$$

The number of states with energy $\leq E$ is then: $N=\frac{V}{3 \pi^{2}}\left(\frac{2 m E}{\hbar^{2}}\right)^{3 / 2}$.
and the density of states becomes: $D(E)=\frac{d N}{d E}=\frac{V}{2 \pi^{2}}\left(\frac{2 m}{\hbar^{2}}\right)^{3 / 2} \sqrt{E}=\frac{3 N}{2 E}$.
The heat capacity of the electrons is approximately 0.01 times the classical expected value $\frac{3}{2} N k$. This is caused by the Pauli exclusion principle and the Fermi-Dirac distribution: only electrons within an energy range $\sim k T$ of the Fermi level are excited thermally. There is a fraction $\approx T / T_{\mathrm{F}}$ excited thermally. The internal energy then becomes:

$$
U \approx N k T \frac{T}{T_{\mathrm{F}}} \text { and } C=\frac{\partial U}{\partial T} \approx N k \frac{T}{T_{\mathrm{F}}}
$$

A more accurate analysis gives: $C_{\text {electrons }}=\frac{1}{2} \pi^{2} N k T / T_{\mathrm{F}} \sim T$. Together with the $T^{3}$ dependence of the thermal heat capacity of the phonons the total thermal heat capacity of metals is described by: $C=\gamma T+A T^{3}$.

### 12.5.2 Electric conductance

The equation of motion for the charge carriers is: $\vec{F}=m d \vec{v} / d t=\hbar d \vec{k} / d t$. The variation of $\vec{k}$ is given by $\delta \vec{k}=\vec{k}(t)-\vec{k}(0)=-e \vec{E} t / \hbar$. If $\tau$ is the characteristic collision time of the electrons, $\delta \vec{k}$ remains stable if $t=\tau$. Then holds: $\langle\vec{v}\rangle=\mu \vec{E}$, with $\mu=e \tau / m$ the mobility of the electrons.
The current in a conductor is given by: $\vec{J}=n q \vec{v}=\sigma \vec{E}=\vec{E} / \rho=n e \mu \vec{E}$. Because for the collision time holds: $1 / \tau=1 / \tau_{L}+1 / \tau_{i}$, where $\tau_{L}$ is the collision time with the lattice phonons and $\tau_{i}$ the collision time with the impurities follows for the resistivity $\rho=\rho_{L}+\rho_{i}$, with $\lim _{T \rightarrow 0} \rho_{L}=0$.

### 12.5.3 The Hall-effect

If a magnetic field is applied $\perp$ to the direction of the current the charge carriers will be pushed aside by the Lorentz force. This results in a magnetic field $\perp$ to the flow direction of the current. If $\vec{J}=J \vec{e}_{x}$ and $\vec{B}=B \vec{e}_{z}$ than $E_{y} / E_{x}=\mu B$. The Hall coefficient is defined by: $R_{\mathrm{H}}=E_{y} / J_{x} B$, and $R_{\mathrm{H}}=-1 / n e$ if $J_{x}=n e \mu E_{x}$. The Hall voltage is given by: $V_{\mathrm{H}}=B v b=I B / n e h$ where $b$ is the width of the material and $h$ de height.

### 12.5.4 Thermal heat conductivity

With $\ell=v_{F} \tau$ the mean free path of the electrons follows from $\kappa=\frac{1}{3} C\langle v\rangle \ell: \kappa_{\text {electrons }}=\pi^{2} n k^{2} T \tau / 3 m$. From this follows for the Wiedemann-Franz ratio: $\kappa / \sigma=\frac{1}{3}(\pi k / e)^{2} T$.

### 12.6 Energy bands

In the tight-bond approximation it is assumed that $\psi=\mathrm{e}^{i k n a} \phi(x-n a)$. From this follows for the energy: $\langle E\rangle=\langle\psi| H|\psi\rangle=E_{\text {at }}-\alpha-2 \beta \cos (k a)$. So this gives a cosine superimposed on the atomic energy, which can often be approximated by a harmonic oscillator. If it is assumed that the electron is nearly free one can postulate: $\psi=\exp (i \vec{k} \cdot \vec{r})$. This is a traveling wave. This wave can be decomposed into two standing waves:

$$
\begin{aligned}
& \psi(+)=\exp (i \pi x / a)+\exp (-i \pi x / a)=2 \cos (\pi x / a) \\
& \psi(-)=\exp (i \pi x / a)-\exp (-i \pi x / a)=2 i \sin (\pi x / a)
\end{aligned}
$$

The probability density $|\psi(+)|^{2}$ is high near the atoms of the lattice and low in between. The probability density $|\psi(-)|^{2}$ is low near the atoms of the lattice and high in between. Hence the energy of $\psi(+)$ is also lower than the energy of $\psi)(-)$. Suppose that $U(x)=U \cos (2 \pi x / a)$, than the bandgap is given by:

$$
E_{\text {gap }}=\int_{0}^{1} U(x)\left[|\psi(+)|^{2}-|\psi(-)|^{2}\right] d x=U
$$

### 12.7 Semiconductors

The band structures and the transitions between them of direct and indirect semiconductors are shown in the figures below. Here it is assumed that the momentum of the absorbed photon can be neglected. For an indirect semiconductor a transition from the valence- to the conduction band is also possible if the energy of the absorbed photon is smaller than the band gap: then, also a phonon is absorbed.


This difference can also be observed in the absorption spectra:


So indirect semiconductors, like Si and Ge , cannot emit any light and are therefore not usable to fabricate lasers. When light is absorbed holds: $\vec{k}_{\mathrm{h}}=-\vec{k}_{\mathrm{e}}, E_{\mathrm{h}}\left(\vec{k}_{\mathrm{h}}\right)=-E_{\mathrm{e}}\left(\vec{k}_{\mathrm{e}}\right), \vec{v}_{\mathrm{h}}=\vec{v}_{\mathrm{e}}$ and $m_{\mathrm{h}}=-m_{\mathrm{e}}^{*}$ if the conduction band and the valence band have the same structure.

Instead of the normal electron mass one has to use the effective mass within a lattice. It is defined by:

$$
m^{*}=\frac{F}{a}=\frac{d p / d t}{d v_{\mathrm{g}} / d t}=\hbar \frac{d K}{d v_{\mathrm{g}}}=\hbar^{2}\left(\frac{d^{2} E}{d k^{2}}\right)^{-1}
$$

with $E=\hbar \omega$ and $v_{\mathrm{g}}=d \omega / d k$ and $p=\hbar k$.
With the distribution function $f_{\mathrm{e}}(E) \approx \exp ((\mu-E) / k T)$ for the electrons and $f_{\mathrm{h}}(E)=1-f_{\mathrm{e}}(E)$ for the holes the density of states is given by:

$$
D(E)=\frac{1}{2 \pi^{2}}\left(\frac{2 m^{*}}{\hbar^{2}}\right)^{3 / 2} \sqrt{E-E_{c}}
$$

with $E_{c}$ the energy at the edge of the conductance band. From this follows for the concentrations of the holes $p$ and the electrons $n$ :

$$
n=\int_{E_{c}}^{\infty} D_{\mathrm{e}}(E) f_{\mathrm{e}}(E) d E=2\left(\frac{m^{*} k T}{2 \pi \hbar^{2}}\right)^{3 / 2} \exp \left(\frac{\mu-E_{c}}{k T}\right)
$$

For the product $n p$ follows: $n p=4\left(\frac{k T}{2 \pi \hbar^{2}}\right)^{3} \sqrt{m_{\mathrm{e}}^{*} m_{\mathrm{h}}} \exp \left(-\frac{E_{\mathrm{g}}}{k T}\right)$
For an intrinsic (no impurities) semiconductor holds: $n_{\mathrm{i}}=p_{\mathrm{i}}$, for a $n-$ type holds: $n>p$ and in a $p-t y p e$ holds: $n<p$.
An exciton is a bound electron-hole pair, rotating on each other as in positronium. The excitation energy of an exciton is smaller than the bandgap because the energy of an exciton is lower than the energy of a free electron and a free hole. This causes a peak in the absorption just under $E_{\mathrm{g}}$.

### 12.8 Superconductivity

### 12.8.1 Description

A superconductor is characterized by a zero resistivity if certain quantities are smaller than some critical values: $T<T_{\mathrm{c}}, I<I_{\mathrm{c}}$ and $H<H_{\mathrm{c}}$. The BCS-model predicts for the transition temperature $T_{\mathrm{c}}$ :

$$
T_{\mathrm{c}}=1.14 \Theta_{\mathrm{D}} \exp \left(\frac{-1}{U D\left(E_{\mathrm{F}}\right)}\right)
$$

while experiments find for $H_{\mathrm{c}}$ approximately:

$$
H_{\mathrm{c}}(T) \approx H_{\mathrm{c}}\left(T_{\mathrm{c}}\right)\left(1-\frac{T^{2}}{T_{\mathrm{c}}^{2}}\right)
$$

Within a superconductor the magnetic field is 0 : the Meissner effect.
There are type I and type II superconductors. Because the Meissner effect implies that a superconductor is a perfect diamagnet holds in the superconducting state: $\vec{H}=\mu_{0} \vec{M}$. This holds for a type I superconductor, for a type II superconductor this only holds to a certain value $H_{\mathrm{c} 1}$, for higher values of $H$ the superconductor is in a vortex state to a value $H_{\mathrm{c} 2}$, which can be 100 times $H_{\mathrm{c} 1}$. If $H$ becomes larger than $H_{\mathrm{c} 2}$ the superconductor becomes a normal conductor. This is shown in the figures below.



The transition to a superconducting state is a second order thermodynamic state transition. This means that there is a twist in the $T-S$ diagram and a discontinuity in the $C_{X}-T$ diagram.

### 12.8.2 The Josephson effect

For the Josephson effect one considers two superconductors, separated by an insulator. The electron wavefunction in one superconductor is $\psi_{1}$, in the other $\psi_{2}$. The Schrödinger equations in both superconductors is set equal:

$$
i \hbar \frac{\partial \psi_{1}}{\partial t}=\hbar T \psi_{2}, i \hbar \frac{\partial \psi_{2}}{\partial t}=\hbar T \psi_{1}
$$

$\hbar T$ is the effect of the coupling of the electrons, or the transfer interaction through the insulator. The electron wavefunctions are written as $\psi_{1}=\sqrt{n_{1}} \exp \left(i \theta_{1}\right)$ and $\psi_{2}=\sqrt{n_{2}} \exp \left(i \theta_{2}\right)$. Because a Cooper pair exist of two electrons holds: $\psi \sim \sqrt{n}$. From this follows, if $n_{1} \approx n_{2}$ :

$$
\frac{\partial \theta_{1}}{\partial t}=\frac{\partial \theta_{2}}{\partial t} \text { and } \frac{\partial n_{2}}{\partial t}=-\frac{\partial n_{1}}{\partial t}
$$

The Josephson effect results in a current density through the insulator depending on the phase difference as: $J=J_{0} \sin \left(\theta_{2}-\theta_{1}\right)=J_{0} \sin (\delta)$, where $J_{0} \sim T$. With an AC-voltage across the junction the Schrödinger equations become:

$$
i \hbar \frac{\partial \psi_{1}}{\partial t}=\hbar T \psi_{2}-e V \psi_{1} \text { and } i \hbar \frac{\partial \psi_{2}}{\partial t}=\hbar T \psi_{1}+e V \psi_{2}
$$

This gives: $J=J_{0} \sin \left(\theta_{2}-\theta_{1}-\frac{2 e V t}{\hbar}\right)$.
Hence there is an oscillation with $\omega=2 e V / \hbar$.

### 12.8.3 Flux quantisation in a superconducting ring

For the current density in general holds: $\vec{J}=q \psi^{*} \vec{v} \psi=\frac{n q}{m}[\hbar \vec{\nabla} \theta-q \vec{A}]$
From the Meissner effect, $\vec{B}=0$ and $\vec{J}=0$, follows: $\hbar \vec{\nabla} \theta=q \vec{A} \Rightarrow \oint \vec{\nabla} \theta d l=\theta_{2}-\theta_{1}=2 \pi s$ with $s \in I N$. Because: $\oint \vec{A} d l=\iint(\operatorname{rot} \vec{A}, \vec{n}) d \sigma=\iint(\vec{B}, \vec{n}) d \sigma=\Psi$ follows: $\Psi=2 \pi \hbar s / q$. The size of a flux quantum follows by setting $s=1: \Psi=2 \pi \hbar / e=2.0678 \cdot 10^{-15} \mathrm{Tm}^{2}$.

### 12.8.4 Macroscopic quantum interference

From $\theta_{2}-\theta_{1}=2 e \Psi / \hbar$ follows for two parallel junctions: $\delta_{b}-\delta_{\mathrm{a}}=\frac{2 e \Psi}{\hbar}$, so
$J=J_{\mathrm{a}}+J_{\mathrm{b}}=2 J_{0} \sin \left(\delta_{0} \cos \left(\frac{e \Psi}{\hbar}\right)\right)$ This gives maxima if $e \Psi / \hbar=s \pi$.

### 12.8.5 The London equation

A current density in a superconductor proportional to the vector potential $\vec{A}$ is postulated:

$$
\vec{J}=\frac{-\vec{A}}{\mu_{0} \lambda_{L}^{2}} \text { or } \operatorname{rot} \vec{J}=\frac{-\vec{B}}{\mu_{0} \lambda_{L}^{2}}
$$

where $\lambda_{L}=\sqrt{\varepsilon_{0} m c^{2} / n q^{2}}$. From this follows: $\nabla^{2} \vec{B}=\vec{B} / \lambda_{L}^{2}$.
The Meissner effect is the solution of this equation: $\vec{B}(x)=B_{0} \exp \left(-x / \lambda_{L}\right)$. Magnetic fields within a superconductor drop exponentially.

### 12.8.6 The BCS model

The BCS model can explain superconductivity in metals. (So far there is no explanation for high- $T_{\mathrm{c}}$ superconductance).

A new ground state where the electrons behave like independent fermions is postulated. Because of the interaction with the lattice these pseudo-particles exhibit a mutual attraction. This causes two electrons with opposite spin to combine to a Cooper pair. It can be proved that this ground state is perfect diamagnetic.

The infinite conductivity is more difficult to explain because a ring with a persisting current is not a real equilibrium: a state with zero current has a lower energy. Flux quantization prevents transitions between these states. Flux quantization is related to the existence of a coherent many-particle wavefunction. A flux quantum is the equivalent of about $10^{4}$ electrons. So if the flux has to change with one flux quantum there has to occur a transition of many electrons, which is very improbable, or the system must go through intermediary states where the flux is not quantized so they have a higher energy. This is also very improbable.
Some useful mathematical relations are:

$$
\int_{0}^{\infty} \frac{x d x}{\mathrm{e}^{a x}+1}=\frac{\pi^{2}}{12 a^{2}}, \int_{-\infty}^{\infty} \frac{x^{2} d x}{\left(\mathrm{e}^{x}+1\right)^{2}}=\frac{\pi^{2}}{3}, \quad \int_{0}^{\infty} \frac{x^{3} d x}{\mathrm{e}^{x}+1}=\frac{\pi^{4}}{15}
$$

And, when $\sum_{n=0}^{\infty}(-1)^{n}=\frac{1}{2}$ follows: $\int_{0}^{\infty} \sin (p x) d x=\int_{0}^{\infty} \cos (p x) d x=\frac{1}{p}$.

## Chapter 13

## Theory of groups

### 13.1 Introduction

### 13.1.1 Definition of a group

$\mathcal{G}$ is a group for the operation $\bullet$ if:

1. $\forall_{A, B \in \mathcal{G}} \Rightarrow A \bullet B \in \mathcal{G}: \mathcal{G}$ is closed.
2. $\forall_{A, B, C \in \mathcal{G}} \Rightarrow(A \bullet B) \bullet C=A \bullet(B \bullet C)$ : $\mathcal{G}$ obeys the associative law.
3. $\exists_{E \in \mathcal{G}}$ so that $\forall_{A \in \mathcal{G}} A \bullet E=E \bullet A=A$ : $\mathcal{G}$ has a unit element.
4. $\forall_{A \in \mathcal{G}} \exists_{A^{-1} \in \mathcal{G}}$ so that $A \bullet A^{-1}=E$ : Each element in $\mathcal{G}$ has an inverse.

If also holds:
5. $\forall_{A, B \in \mathcal{G}} \Rightarrow A \bullet B=B \bullet A$ the group is called Abelian or commutative.

### 13.1.2 The Cayley table

Each element arises only once in each row and column of the Cayley- or multiplication table: because $E A_{i}=$ $A_{k}^{-1}\left(A_{k} A_{i}\right)=A_{i}$ each $A_{i}$ appears once. There are $h$ positions in each row and column when there are $h$ elements in the group so each element appears only once.

### 13.1.3 Conjugated elements, subgroups and classes

$B$ is conjugate to $A$ if $\exists_{X \in \mathcal{G}}$ such that $B=X A X^{-1}$. Then $A$ is also conjugate to $B$ because $B=$ $\left(X^{-1}\right) A\left(X^{-1}\right)^{-1}$.
If $B$ and $C$ are conjugate to $A, B$ is also conjugate with $C$.
A subgroup is a subset of $\mathcal{G}$ which is also a group w.r.t. the same operation.
A conjugacy class is the maximum collection of conjugated elements. Each group can be split up in conjugacy classes. Some theorems:

- All classes are completely disjoint.
- $E$ is a class itself: for each other element in this class would hold: $A=X E X^{-1}=E$.
- $E$ is the only class which is also a subgroup because all other classes have no unit element.
- In an Abelian group each element is a separate class.

The physical interpretation of classes: elements of a group are usually symmetry operations which map a symmetrical object into itself. Elements of one class are then the same kind of operations. The opposite need not to be true.

### 13.1.4 Isomorfism and homomorfism; representations

Two groups are isomorphic if they have the same multiplication table. The mapping from group $\mathcal{G}_{1}$ to $\mathcal{G}_{2}$, so that the multiplication table remains the same is a homomorphic mapping. It need not be isomorphic.
A representation is a homomorphic mapping of a group to a group of square matrices with the usual matrix multiplication as the combining operation. This is symbolized by $\Gamma$. The following holds:

$$
\Gamma(E)=\mathbb{I}, \quad \Gamma(A B)=\Gamma(A) \Gamma(B), \Gamma\left(A^{-1}\right)=[\Gamma(A)]^{-1}
$$

For each group there are 3 possibilities for a representation:

1. A faithful representation: all matrices are different.
2. The representation $A \rightarrow \operatorname{det}(\Gamma(A))$.
3. The identical representation: $A \rightarrow 1$.

An equivalent representation is obtained by performing an unitary base transformation: $\Gamma^{\prime}(A)=S^{-1} \Gamma(A) S$.

### 13.1.5 Reducible and irreducible representations

If the same unitary transformation can bring all matrices of a representation $\Gamma$ in the same block structure the representation is called reducible:

$$
\Gamma(A)=\left(\begin{array}{cc}
\Gamma^{(1)}(A) & 0 \\
0 & \Gamma^{(2)}(A)
\end{array}\right)
$$

This is written as: $\Gamma=\Gamma^{(1)} \oplus \Gamma^{(2)}$. If this is not possible the representation is called irreducible.
The number of irreducible representations equals the number of conjugacy classes.

### 13.2 The fundamental orthogonality theorem

### 13.2.1 Schur's lemma

Lemma: Each matrix which commutes with all matrices of an irreducible representation is a constant $\times I I$, where $I I$ is the unit matrix. The opposite is (of course) also true.
Lemma: If there exists a matrix $M$ so that for two irreducible representations of group $\mathcal{G}, \gamma^{(1)}\left(A_{i}\right)$ and $\overline{\gamma^{(2)}\left(A_{i}\right)}$, holds: $M \gamma^{(1)}\left(A_{i}\right)=\gamma^{(2)}\left(A_{i}\right) M$, than the representations are equivalent, or $M=\underline{\underline{0}}$.

### 13.2.2 The fundamental orthogonality theorem

For a set of unequivalent, irreducible, unitary representations holds that, if $h$ is the number of elements in the group and $\ell_{i}$ is the dimension of the $i$ th representation:

$$
\sum_{R \in \mathcal{G}} \Gamma_{\mu \nu}^{(i) *}(R) \Gamma_{\alpha \beta}^{(j)}(R)=\frac{h}{\ell_{i}} \delta_{i j} \delta_{\mu \alpha} \delta_{\nu \beta}
$$

### 13.2.3 Character

The character of a representation is given by the trace of the matrix and is therefore invariant for base transformations: $\chi^{(j)}(R)=\operatorname{Tr}\left(\Gamma^{(j)}(R)\right)$
Also holds, with $N_{k}$ the number of elements in a conjugacy class: $\sum_{k} \chi^{(i) *}\left(C_{k}\right) \chi^{(j)}\left(C_{k}\right) N_{k}=h \delta_{i j}$
Theorem: $\sum_{i=1}^{n} \ell_{i}^{2}=h$

### 13.3 The relation with quantum mechanics

### 13.3.1 Representations, energy levels and degeneracy

Consider a set of symmetry transformations $\vec{x}^{\prime}=R \vec{x}$ which leave the Hamiltonian $\mathcal{H}$ invariant. These transformations are a group. An isomorfic operation on the wavefunction is given by: $P_{R} \psi(\vec{x})=\psi\left(R^{-1} \vec{x}\right)$. This is considered an active rotation. These operators commute with $\mathcal{H}: P_{R} \mathcal{H}=\mathcal{H} P_{R}$, and leave the volume element unchanged: $d(R \vec{x})=d \vec{x}$.
$P_{R}$ is the symmetry group of the physical system. It causes degeneracy: if $\psi_{n}$ is a solution of $\mathcal{H} \psi_{n}=E_{n} \psi_{n}$ than also holds: $\mathcal{H}\left(P_{R} \psi_{n}\right)=E_{n}\left(P_{R} \psi_{n}\right)$. A degeneracy which is not the result of a symmetry is called an accidental degeneracy.
Assume an $\ell_{n}$-fold degeneracy at $E_{n}$ : then choose an orthonormal set $\psi_{\nu}^{(n)}, \nu=1,2, \ldots, \ell_{n}$. The function $P_{R} \psi_{\nu}^{(n)}$ is in the same subspace: $P_{R} \psi_{\nu}^{(n)}=\sum_{\kappa=1}^{\ell_{n}} \psi_{\kappa}^{(n)} \Gamma_{\kappa \nu}^{(n)}(R)$
where $\Gamma^{(n)}$ is an irreducible, unitary representation of the symmetry group $\mathcal{G}$ of the system. Each $n$ corresponds with another energy level. One can purely mathematical derive irreducible representations of a symmetry group and label the energy levels with a quantum number this way. A fixed choice of $\Gamma^{(n)}(R)$ defines the base functions $\psi_{\nu}^{(n)}$. This way one can also label each separate base function with a quantum number.
Particle in a periodical potential: the symmetry operation is a cyclic group: note the operator describing one translation over one unit as $A$. Then: $\mathcal{G}=\left\{A, A^{2}, A^{3}, \ldots, A^{h}=E\right\}$.
The group is Abelian so all irreducible representations are one-dimensional. For $0 \leq p \leq h-1$ follows:

$$
\Gamma^{(p)}\left(A^{n}\right)=\mathrm{e}^{2 \pi i p n / h}
$$

If one defines: $k=-\frac{2 \pi p}{a h}\left(\bmod \frac{2 \pi}{a}\right)$, so: $P_{A} \psi_{p}(x)=\psi_{p}(x-a)=\mathrm{e}^{2 \pi i p / h} \psi_{p}(x)$, this gives Bloch's theorem: $\psi_{k}(x)=u_{k}(x) \mathrm{e}^{i k x}$, with $u_{k}(x \pm a)=u_{k}(x)$.

### 13.3.2 Breaking of degeneracy by a perturbation

Suppose the unperturbed system has Hamiltonian $\mathcal{H}_{0}$ and symmetry group $\mathcal{G}_{0}$. The perturbed system has $\mathcal{H}=\mathcal{H}_{0}+\mathcal{V}$, and symmetry group $\mathcal{G} \subset \mathcal{G}_{0}$. If $\Gamma^{(n)}(R)$ is an irreducible representation of $\mathcal{G}_{0}$, it is also a representation of $\mathcal{G}$ but not all elements of $\Gamma^{(n)}$ in $\mathcal{G}_{0}$ are also in $\mathcal{G}$. The representation then usually becomes reducible: $\Gamma^{(n)}=\Gamma^{\left(n_{1}\right)} \oplus \Gamma^{\left(n_{2}\right)} \oplus \ldots$. The degeneracy is then (possibly partially) removed: see the figure below.


Spectrum $\mathcal{H}_{0}$
Spectrum $\mathcal{H}$
Theorem: The set of $\ell_{n}$ degenerated eigenfunctions $\psi_{\nu}^{(n)}$ with energy $E_{n}$ is a basis for an $\ell_{n}$-dimensional irreducible representation $\Gamma^{(n)}$ of the symmetry group.

### 13.3.3 The construction of a base function

Each function $F$ in configuration space can be decomposed into symmetry types: $F=\sum_{j=1}^{n} \sum_{\kappa=1}^{\ell_{j}} f_{\kappa}^{(j)}$
The following operator extracts the symmetry types:

$$
\left(\frac{\ell_{j}}{h} \sum_{R \in \mathcal{G}} \Gamma_{\kappa \kappa}^{(j) *}(R) P_{R}\right) F=f_{\kappa}^{(j)}
$$

This is expressed as: $f_{\kappa}^{(j)}$ is the part of $F$ that transforms according to the $\kappa^{\text {th }}$ row of $\Gamma^{(j)}$.
$F$ can also be expressed in base functions $\varphi: F=\sum_{a j \kappa} c_{a j \kappa} \varphi_{\kappa}^{(a j)}$. The functions $f_{\kappa}^{(j)}$ are in general not transformed into each other by elements of the group. However, this does happen if $c_{j a \kappa}=c_{j a}$.

Theorem: Two wavefunctions transforming according to non-equivalent unitary representations or according to different rows of an unitary irreducible representation are orthogonal:
$\left\langle\varphi_{\kappa}^{(i)} \mid \psi_{\lambda}^{(j)}\right\rangle \sim \delta_{i j} \delta_{\kappa \lambda}$, and $\left\langle\varphi_{\kappa}^{(i)} \mid \psi_{\kappa}^{(i)}\right\rangle$ is independent of $\kappa$.

### 13.3.4 The direct product of representations

Consider a physical system existing of two subsystems. The subspace $D^{(i)}$ of the system transforms according to $\Gamma^{(i)}$. Basefunctions are $\varphi_{\kappa}^{(i)}\left(\vec{x}_{i}\right), 1 \leq \kappa \leq \ell_{i}$. Now form all $\ell_{1} \times \ell_{2}$ products $\varphi_{\kappa}^{(1)}\left(\vec{x}_{1}\right) \varphi_{\lambda}^{(2)}\left(\vec{x}_{2}\right)$. These define a space $D^{(1)} \otimes D^{(2)}$.

These product functions transform as:

$$
P_{R}\left(\varphi_{\kappa}^{(1)}\left(\vec{x}_{1}\right) \varphi_{\lambda}^{(2)}\left(\vec{x}_{2}\right)\right)=\left(P_{R} \varphi_{\kappa}^{(1)}\left(\vec{x}_{1}\right)\right)\left(P_{R} \varphi_{\lambda}^{(2)}\left(\vec{x}_{2}\right)\right)
$$

In general the space $D^{(1)} \otimes D^{(2)}$ can be split up in a number of invariant subspaces:

$$
\Gamma^{(1)} \otimes \Gamma^{(2)}=\sum_{i} n_{i} \Gamma^{(i)}
$$

A useful tool for this reduction is that for the characters hold:

$$
\chi^{(1)}(R) \chi^{(2)}(R)=\sum_{i} n_{i} \chi^{(i)}(R)
$$

### 13.3.5 Clebsch-Gordan coefficients

With the reduction of the direct-product matrix w.r.t. the $\operatorname{basis} \varphi_{\kappa}^{(i)} \varphi_{\lambda}^{(j)}$ one uses a new basis $\varphi_{\mu}^{(a \kappa)}$. These base functions lie in subspaces $D^{(a k)}$. The unitary base transformation is given by:

$$
\varphi_{\mu}^{(a k)}=\sum_{\kappa \lambda} \varphi_{\kappa}^{(i)} \varphi_{\lambda}^{(j)}(i \kappa j \lambda \mid a k \mu)
$$

and the inverse transformation by: $\varphi_{\kappa}^{(i)} \varphi_{\lambda}^{(j)}=\sum_{a k \mu} \varphi_{\mu}^{(a \kappa)}(a k \mu \mid i \kappa j \lambda)$
In essence the Clebsch-Gordan coefficients are dot products: $(i \kappa j \lambda \mid a k \mu):=\left\langle\varphi_{k}^{(i)} \varphi_{\lambda}^{(j)} \mid \varphi_{\mu}^{(a k)}\right\rangle$

### 13.3.6 Symmetric transformations of operators, irreducible tensor operators

Observables (operators) transform as follows under symmetry transformations: $A^{\prime}=P_{R} A P_{R}^{-1}$. If a set of operators $A_{\kappa}^{(j)}$ with $0 \leq \kappa \leq \ell_{j}$ transform into each other under the transformations of $\mathcal{G}$ holds:

$$
P_{R} A_{\kappa}^{(j)} P_{R}^{-1}=\sum_{\nu} A_{\nu}^{(j)} \Gamma_{\nu \kappa}^{(j)}(R)
$$

If $\Gamma^{(j)}$ is irreducible they are called irreducible tensor operators $A^{(j)}$ with components $A_{\kappa}^{(j)}$.
An operator can also be decomposed into symmetry types: $A=\sum_{j k} a_{k}^{(j)}$, with:

$$
a_{\kappa}^{(j)}=\left(\frac{\ell_{j}}{h} \sum_{R \in \mathcal{G}} \Gamma_{\kappa \kappa}^{(j) *}(R)\right)\left(P_{R} A P_{R}^{-1}\right)
$$

Theorem: Matrix elements $H_{i j}$ of the operator $\mathcal{H}$ which is invariant under $\forall_{A \in \mathcal{G}}$, are 0 between states which transform according to non-equivalent irreducible unitary representations or according to different rows of such a representation. Further $\left\langle\varphi_{\kappa}^{(i)}\right| \mathcal{H}\left|\psi_{\kappa}^{(i)}\right\rangle$ is independent of $\kappa$. For $\mathcal{H}=1$ this becomes the previous theorem.
This is applied in quantum mechanics in perturbation theory and variational calculus. Here one tries to diagonalize $\mathcal{H}$. Solutions can be found within each category of functions $\varphi_{\kappa}^{(i)}$ with common $i$ and $\kappa$ : $\mathcal{H}$ is already diagonal in categories as a whole.
Perturbation calculus can be applied independent within each category. With variational calculus the try function can be chosen within a separate category because the exact eigenfunctions transform according to a row of an irreducible representation.

### 13.3.7 The Wigner-Eckart theorem

Theorem: The matrix element $\left\langle\varphi_{\lambda}^{(i)}\right| A_{\kappa}^{(j)}\left|\psi_{\mu}^{(k)}\right\rangle$ can only be $\neq 0$ if $\Gamma^{(j)} \otimes \Gamma^{(k)}=\ldots \oplus \Gamma^{(i)} \oplus \ldots$. If this is the case holds (if $\Gamma^{(i)}$ appears only once, otherwise one has to sum over $a$ ):

$$
\left\langle\varphi_{\lambda}^{(i)}\right| A_{\kappa}^{(j)}\left|\psi_{\mu}^{(k)}\right\rangle=(i \lambda \mid j \kappa k \mu)\left\langle\varphi^{(i)}\left\|A^{(j)}\right\| \psi^{(k)}\right\rangle
$$

This theorem can be used to determine selection rules: the probability of a dipole transition is given by ( $\vec{\epsilon}$ is the direction of polarization of the radiation):

$$
P_{D}=\frac{8 \pi^{2} e^{2} f^{3}\left|r_{12}\right|^{2}}{3 \hbar \varepsilon_{0} c^{3}} \text { with } r_{12}=\left\langle l_{2} m_{2}\right| \vec{\epsilon} \cdot \vec{r}\left|l_{1} m_{1}\right\rangle
$$

Further it can be used to determine intensity ratios: if there is only one value of $a$ the ratio of the matrix elements are the Clebsch-Gordan coefficients. For more $a$-values relations between the intensity ratios can be stated. However, the intensity ratios are also dependent on the occupation of the atomic energy levels.

### 13.4 Continuous groups

Continuous groups have $h=\infty$. However, not all groups with $h=\infty$ are continuous, e.g. the translation group of an spatially infinite periodic potential is not continuous but does have $h=\infty$.

### 13.4.1 The 3-dimensional translation group

For the translation of wavefunctions over a distance $a$ holds: $P_{a} \psi(x)=\psi(x-a)$. Taylor expansion near $x$ gives:

$$
\psi(x-a)=\psi(x)-a \frac{d \psi(x)}{d x}+\frac{1}{2} a^{2} \frac{d^{2} \psi(x)}{d x^{2}}-+\ldots
$$

Because the momentum operator in quantum mechanics is given by: $p_{x}=\frac{\hbar}{i} \frac{\partial}{\partial x}$, this can be written as:

$$
\psi(x-a)=\mathrm{e}^{-i a p_{x} / \hbar} \psi(x)
$$

### 13.4.2 The 3-dimensional rotation group

This group is called $\mathrm{SO}(3)$ because a faithful representation can be constructed from orthogonal $3 \times 3$ matrices with a determinant of +1 .
For an infinitesimal rotation around the $x$-axis holds:

$$
\begin{aligned}
P_{\delta \theta_{x}} \psi(x, y, z) & \approx \psi\left(x, y+z \delta \theta_{x}, z-y \delta \theta_{x}\right) \\
& =\psi(x, y, z)+\left(z \delta \theta_{x} \frac{\partial}{\partial y}-y \delta \theta_{x} \frac{\partial}{\partial z}\right) \psi(x, y, z) \\
& =\left(1-\frac{i \delta \theta_{x} L_{x}}{\hbar}\right) \psi(x, y, z)
\end{aligned}
$$

Because the angular momentum operator is given by: $L_{x}=\frac{\hbar}{i}\left(z \frac{\partial}{\partial y}-y \frac{\partial}{\partial z}\right)$.
So in an arbitrary direction holds: $\quad$ Rotations: $\quad P_{\alpha, \vec{n}}=\exp (-i \alpha(\vec{n} \cdot \vec{J}) / \hbar)$
Translations: $\quad P_{a, \vec{n}}=\exp (-i a(\vec{n} \cdot \vec{p}) / \hbar)$
$J_{x}, J_{y}$ and $J_{z}$ are called the generators of the 3-dim. rotation group, $p_{x}, p_{y}$ and $p_{z}$ are called the generators of the 3-dim. translation group.

The commutation rules for the generators can be derived from the properties of the group for multiplications: translations are interchangeable $\leftrightarrow p_{x} p_{y}-p_{y} p_{x}=0$.
Rotations are not generally interchangeable: consider a rotation around axis $\vec{n}$ in the $x z$-plane over an angle $\alpha$. Then holds: $P_{\alpha, \vec{n}}=P_{-\theta, y} P_{\alpha, x} P_{\theta, y}$, so:

$$
\mathrm{e}^{-i \alpha(\vec{n} \cdot \vec{J}) / \hbar}=\mathrm{e}^{i \theta J_{y} / \hbar} \mathrm{e}^{-i \alpha J_{x} / \hbar} \mathrm{e}^{-i \theta J_{y} / \hbar}
$$

If $\alpha$ and $\theta$ are very small and are expanded to second order, and the corresponding terms are put equal with $\vec{n} \cdot \vec{J}=J_{x} \cos \theta+J_{z} \sin \theta$, it follows from the $\alpha \theta$ term: $J_{x} J_{y}-J_{y} J_{x}=i \hbar J_{z}$.

### 13.4.3 Properties of continuous groups

The elements $R\left(p_{1}, \ldots, p_{n}\right)$ depend continuously on parameters $p_{1}, \ldots, p_{n}$. For the translation group this are e.g. $a n_{x}, a n_{y}$ and $a n_{z}$. It is demanded that the multiplication and inverse of an element $R$ depend continuously on the parameters of $R$.
The statement that each element arises only once in each row and column of the Cayley table holds also for continuous groups. The notion conjugacy class for continuous groups is defined equally as for discrete groups. The notion representation is fitted by demanding continuity: each matrix element depends continuously on $p_{i}(R)$.

Summation over all group elements is for continuous groups replaced by an integration. If $f(R)$ is a function defined on $\mathcal{G}$, e.g. $\Gamma_{\alpha \beta}(R)$, holds:

$$
\int_{\mathcal{G}} f(R) d R:=\int_{p_{1}} \cdots \int_{p_{n}} f\left(R\left(p_{1}, \ldots, p_{n}\right)\right) g\left(R\left(p_{1}, \ldots, p_{n}\right)\right) d p_{1} \cdots d p_{n}
$$

Here, $g(R)$ is the density function.
Because of the properties of the Cayley table is demanded: $\int f(R) d R=\int f(S R) d R$. This fixes $g(R)$ except for a constant factor. Define new variables $p^{\prime}$ by: $S R\left(p_{i}\right)=R\left(p_{i}^{\prime}\right)$. If one writes: $d V:=d p_{1} \cdots d p_{n}$ holds:

$$
g(S)=g(E) \frac{d V}{d V^{\prime}}
$$

Here, $\frac{d V}{d V^{\prime}}$ is the Jacobian: $\frac{d V}{d V^{\prime}}=\operatorname{det}\left(\frac{\partial p_{i}}{\partial p_{j}^{\prime}}\right)$, and $g(E)$ is constant.
For the translation group holds: $g(\vec{a})=$ constant $=g(\overrightarrow{0})$ because $g(a \vec{n}) d \vec{a}^{\prime}=g(\overrightarrow{0}) d \vec{a}$ and $d \vec{a}^{\prime}=d \vec{a}$.
This leads to the fundamental orthogonality theorem:

$$
\int_{\mathcal{G}} \Gamma_{\mu \nu}^{(i) *}(R) \Gamma_{\alpha \beta}^{(j)}(R) d R=\frac{1}{\ell_{i}} \delta_{i j} \delta_{\mu \alpha} \delta_{\nu \beta} \int_{\mathcal{G}} d R
$$

and for the characters hold:

$$
\int_{\mathcal{G}} \chi^{(i) *}(R) \chi^{(j)}(R) d R=\delta_{i j} \int_{\mathcal{G}} d R
$$

Compact groups are groups with a finite group volume: $\int_{\mathcal{G}} d R<\infty$.

### 13.5 The group $\mathrm{SO}(3)$

One can take 2 parameters for the direction of the rotational axis and one for the angle of rotation $\varphi$. The parameter space is a collection points $\varphi \vec{n}$ within a sphere with radius $\pi$. The diametrical points on this sphere are equivalent because $R_{\vec{n}, \pi}=R_{\vec{n},-\pi}$.

Another way to define parameters is by means of Eulers angles. If $\alpha, \beta$ and $\gamma$ are the 3 Euler angles, defined as:

1. The spherical angles of axis 3 w.r.t. $x y z$ are $\theta, \varphi:=\beta, \alpha$. Now a rotation around axis 3 remains possible.
2. The spherical angles of the $z$-axis w.r.t. 123 are $\theta, \varphi:=\beta, \pi-\gamma$.
then the rotation of a quantum mechanical system is described by:
$\psi \rightarrow \mathrm{e}^{-i \alpha J_{z} \hbar} \mathrm{e}^{-i \beta J_{y} / \hbar} \mathrm{e}^{-i \gamma J_{z} / \hbar} \psi$. So $P_{R}=\mathrm{e}^{-i \varepsilon(\vec{n} \cdot \vec{J}) / \hbar}$.
All irreducible representations of $\mathrm{SO}(3)$ can be constructed from the behaviour of the spherical harmonics $Y_{l m}(\theta, \varphi)$ with $-l \leq m \leq l$ and for a fixed $l$ :

$$
P_{R} Y_{l m}(\theta, \varphi)=\sum_{m^{\prime}} Y_{l m^{\prime}}(\theta, \varphi) D_{m m^{\prime}}^{(l)}(R)
$$

$D^{(l)}$ is an irreducible representation of dimension $2 l+1$. The character of $D^{(l)}$ is given by:

$$
\chi^{(l)}(\alpha)=\sum_{m=-l}^{l} \mathrm{e}^{i m \alpha}=1+2 \sum_{k=0}^{l} \cos (k \alpha)=\frac{\sin \left(\left[l+\frac{1}{2}\right] \alpha\right)}{\sin \left(\frac{1}{2} \alpha\right)}
$$

In the performed derivation $\alpha$ is the rotational angle around the $z$-axis. This expression is valid for all rotations over an angle $\alpha$ because the classes of $\mathrm{SO}(3)$ are rotations around the same angle around an axis with an arbitrary orientation.
Via the fundamental orthogonality theorem for characters one obtains the following expression for the density function (which is normalized so that $g(0)=1$ ):

$$
g(\alpha)=\frac{\sin ^{2}\left(\frac{1}{2} \alpha\right)}{\left(\frac{1}{2} \alpha\right)^{2}}
$$

With this result one can see that the given representations of $\mathrm{SO}(3)$ are the only ones: the character of another representation $\chi^{\prime}$ would have to be $\perp$ to the already found ones, so $\chi^{\prime}(\alpha) \sin ^{2}\left(\frac{1}{2} \alpha\right)=0 \forall \alpha \Rightarrow \chi^{\prime}(\alpha)=0 \forall \alpha$. This is contradictory because the dimension of the representation is given by $\chi^{\prime}(0)$.
Because fermions have an half-odd integer spin the states $\psi_{s m_{s}}$ with $s=\frac{1}{2}$ and $m_{s}= \pm \frac{1}{2}$ constitute a 2-dim. space which is invariant under rotations. A problem arises for rotations over $2 \pi$ :

$$
\psi_{\frac{1}{2} m_{s}} \rightarrow \mathrm{e}^{-2 \pi i S_{z} / \hbar} \psi_{\frac{1}{2} m_{s}}=\mathrm{e}^{-2 \pi i m_{s}} \psi_{\frac{1}{2} m_{s}}=-\psi_{\frac{1}{2} m_{s}}
$$

However, in $\mathrm{SO}(3)$ holds: $R_{z, 2 \pi}=E$. So here holds $E \rightarrow \pm I$. Because observable quantities can always be written as $\langle\phi \mid \psi\rangle$ or $\langle\phi| A|\psi\rangle$, and are bilinear in the states, they do not change sign if the states do. If only one state changes sign the observable quantities do change.
The existence of these half-odd integer representations is connected with the topological properties of $\mathrm{SO}(3)$ : the group is two-fold coherent through the identification $R_{0}=R_{2 \pi}=E$.

### 13.6 Applications to quantum mechanics

### 13.6.1 Vectormodel for the addition of angular momentum

If two subsystems have angular momentum quantum numbers $j_{1}$ and $j_{2}$ the only possible values for the total angular momentum are $J=j_{1}+j_{2}, j_{1}+j_{2}-1, \ldots,\left|j_{1}-j_{2}\right|$. This can be derived from group theory as follows: from $\chi^{\left(j_{1}\right)}(\alpha) \chi^{\left(j_{2}\right)}(\alpha)=\sum_{J} n_{j} \chi^{(J)}(\alpha)$ follows:

$$
D^{\left(j_{1}\right)} \otimes D^{\left(j_{2}\right)}=D^{\left(j_{1}+j_{2}\right)} \oplus D^{\left(j_{1}+j_{2}-1\right)} \oplus \ldots \oplus D^{\left(\left|j_{1}-j_{2}\right|\right)}
$$

The states can be characterized by quantum numbers in two ways: with $j_{1}, m_{1}, j_{2}, m_{2}$ and with $j_{1}, j_{2}, J, M$. The Clebsch-Gordan coefficients, for $\mathrm{SO}(3)$ called the Wigner coefficients, can be chosen real, so:

$$
\begin{aligned}
& \psi_{j_{1} j_{2} J M}=\sum_{m_{1} m_{2}} \psi_{j_{1} m_{1} j_{2} m_{2}}\left(j_{1} m_{1} j_{2} m_{2} \mid J M\right) \\
& \psi_{j_{1} m_{1} j_{2} m_{2}}=\sum_{J M} \psi_{j_{1} j_{2} J M}\left(j_{1} m_{1} j_{2} m_{2} \mid J M\right)
\end{aligned}
$$

### 13.6.2 Irreducible tensor operators, matrixelements and selection rules

## Some examples of the behaviour of operators under $\mathrm{SO}(3)$

1. Suppose $j=0$ : this gives the identical representation with $\ell_{j}=1$. This state is described by a scalar operator. Because $P_{R} A_{0}^{(0)} P_{R}^{-1}=A_{0}^{(0)}$ this operator is invariant, e.g. the Hamiltonian of a free atom. Then holds: $\left\langle J^{\prime} M^{\prime}\right| \mathcal{H}|J M\rangle \sim \delta_{M M^{\prime}} \delta_{J J^{\prime}}$.
2. A vector operator: $\vec{A}=\left(A_{x}, A_{y}, A_{z}\right)$. The cartesian components of a vector operator transform equally as the cartesian components of $\vec{r}$ by definition. So for rotations around the $z$-axis holds:

$$
D\left(R_{\alpha, z}\right)=\left(\begin{array}{ccc}
\cos \alpha & -\sin \alpha & 0 \\
\sin \alpha & \cos \alpha & 0 \\
0 & 0 & 1
\end{array}\right)
$$

The transformed operator has the same matrix elements w.r.t. $P_{R} \psi$ and $P_{R} \phi$ :
$\left\langle P_{R} \psi\right| P_{R} A_{x} P_{R}^{-1}\left|P_{R} \phi\right\rangle=\langle\psi| A_{x}|\phi\rangle$, and $\chi\left(R_{\alpha, z}\right)=1+2 \cos (\alpha)$. According to the equation for characters this means one can choose base operators which transform like $Y_{1 m}(\theta, \varphi)$. These turn out to be the spherical components:

$$
A_{+1}^{(1)}=-\frac{1}{\sqrt{2}}\left(A_{x}+i A_{y}\right), \quad A_{0}^{(1)}=A_{z}, \quad A_{-1}^{(1)}=\frac{1}{\sqrt{2}}\left(A_{x}-i A_{y}\right)
$$

3. A cartesian tensor of rank 2: $T_{i j}$ is a quantity which transforms under rotations like $U_{i} V_{j}$, where $\vec{U}$ and $\vec{V}$ are vectors. So $T_{i j}$ transforms like $P_{R} T_{i j} P_{R}^{-1}=\sum_{k l} T_{k l} D_{k i}(R) D_{l j}(R)$, so like $D^{(1)} \otimes D^{(1)}=$ $D^{(2)} \oplus D^{(1)} \oplus D^{(0)}$. The 9 components can be split in 3 invariant subspaces with dimension $1\left(D^{(0)}\right)$, $3\left(D^{(1)}\right)$ and $5\left(D^{(2)}\right)$. The new base operators are:
I. $\operatorname{Tr}(\underline{\underline{T}})=T_{x x}+T_{y y}+T_{z z}$. This transforms as the scalar $\vec{U} \cdot \vec{V}$, so as $D^{(0)}$.
II. The 3 antisymmetric components $A_{z}=\frac{1}{2}\left(T_{x y}-T_{y x}\right)$, etc. These transform as the vector $\vec{U} \times \vec{V}$, so as $D^{(1)}$.
III. The 5 independent components of the traceless, symmetric tensor $\underline{\underline{S}}$ : $S_{i j}=\frac{1}{2}\left(T_{i j}+T_{j i}\right)-\frac{1}{3} \delta_{i j} \operatorname{Tr}(\underline{\underline{T}})$. These transform as $D^{(2)}$.

## Selection rules for dipole transitions

Dipole operators transform as $D^{(1)}$ : for an electric dipole transfer is the operator $e \vec{r}$, for a magnetic $e(\vec{L}+$ $2 \vec{S}) / 2 m$.
From the Wigner-Eckart theorem follows: $\left\langle J^{\prime} M^{\prime}\right| A_{\kappa}^{(1)}|J M\rangle=0$ except $D^{\left(J^{\prime}\right)}$ is a part of $D^{(1)} \otimes D^{(J)}=$ $D^{(J+1)} \oplus D^{(J)} \oplus D^{(|J-1|)}$. This means that $J^{\prime} \in\{J+1, J,|J-1|\}: J^{\prime}=J$ or $J^{\prime}=J \pm 1$, except $J^{\prime}=J=0$.

## Landé-equation for the anomalous Zeeman splitting

According to Landés model the interaction between a magnetic moment with an external magnetic field is determined by the projection of $\vec{M}$ on $\vec{J}$ because $\vec{L}$ and $\vec{S}$ precede fast around $\vec{J}$. This can also be understood
from the Wigner-Eckart theorem: from this follows that the matrix elements from all vector operators show a certain proportionality. For an arbitrary operator $\vec{A}$ follows:

$$
\left\langle\alpha j m^{\prime}\right| \vec{A}|\alpha j m\rangle=\frac{\langle\alpha j m| \vec{A} \cdot \vec{J}|\alpha j m\rangle}{j(j+1) \hbar^{2}}\left\langle\alpha j m^{\prime}\right| \vec{J}|\alpha j m\rangle
$$

### 13.7 Applications to particle physics

The physics of a system does not change after performing a transformation $\psi^{\prime}=\mathrm{e}^{i \delta} \psi$ where $\delta$ is a constant. This is a global gauge transformation: the phase of the wavefunction changes everywhere by the same amount.

There exists some freedom in the choice of the potentials $\vec{A}$ and $\phi$ at the same $\vec{E}$ and $\vec{B}$ : gauge transformations of the potentials do not change $\vec{E}$ and $\vec{B}$ (See chapter 2 and 10). The solution $\psi^{\prime}$ of the Schrödinger equation with the transformed potentials is: $\psi^{\prime}=\mathrm{e}^{-i q f(\vec{r}, t)} \psi$.

This is a local gauge transformation: the phase of the wavefunction changes different at each position. The physics of the system does not change if $\vec{A}$ and $\phi$ are also transformed. This is now stated as a guide principle: the "right of existence" of the electromagnetic field is to allow local gauge invariance.

The gauge transformations of the EM-field form a group: $\mathrm{U}(1)$, unitary $1 \times 1$-matrices. The split-off of charge in the exponent is essential: it allows one gauge field for all charged particles, independent of their charge.

This concept is generalized: particles have a "special charge" $Q$. The group elements now are $P_{R}=\exp (-i Q \Theta)$.

Other force fields than the electromagnetic field can also be understood this way. The weak interaction together with the electromagnetic interaction can be described by a force field that transforms according to $\mathrm{U}(1) \otimes \mathrm{SU}(2)$, and consists of the photon and three intermediary vector bosons. The colour force is described by $\operatorname{SU}(3)$, and has a gauge field that exists of 8 types of gluons.

In general the group elements are given by $P_{R}=\exp (-i \vec{T} \cdot \vec{\Theta})$, where $\Theta_{n}$ are real constants and $T_{n}$ operators (generators), like $Q$. The commutation rules are given by $\left[T_{i}, T_{j}\right]=i \sum_{k} c_{i j k} T_{k}$. The $c_{i j k}$ are the structure constants of the group. For $\operatorname{SO}(3)$ these constants are $c_{i j k}=\varepsilon_{i j k}$, here $\varepsilon_{i j k}$ is the complete antisymmetric tensor with $\varepsilon_{123}=+1$.

These constants can be found with the help of group product elements: because $\mathcal{G}$ is closed holds: $\mathrm{e}^{i \vec{\Theta} \cdot \vec{T}} \mathrm{e}^{i \vec{\Theta}^{\prime} \cdot \vec{T}} \mathrm{e}^{-i \vec{\Theta} \cdot \vec{T}} \mathrm{e}^{-i \vec{\Theta}^{\prime} \cdot \vec{T}}=\mathrm{e}^{-i \vec{\Theta}^{\prime \prime} \cdot \vec{T}}$. Taylor expansion and setting equal $\Theta^{n} \Theta^{\prime m}$-terms results in the commutation rules.

The group $S U(2)$ has 3 free parameters: because it is unitary there are 4 real conditions over 4 complex parameters, and the determinant has to be +1 , remaining 3 free parameters.
Each unitary matrix $U$ can be written as: $U=\mathrm{e}^{-i H}$. Here, $H$ is a Hermitian matrix. Further it always holds that: $\operatorname{det}(U)=\mathrm{e}^{-i \operatorname{Tr}(H)}$.
For each matrix of $\operatorname{SU}(2)$ holds that $\operatorname{Tr}(H)=0$. Each Hermitian, traceless $2 \times 2$ matrix can be written as a linear combination of the 3 Pauli-matrices $\sigma_{i}$. So these matrices are a choice for the operators of $\mathrm{SU}(2)$. One can write: $\mathrm{SU}(2)=\left\{\exp \left(-\frac{1}{2} i \vec{\sigma} \cdot \vec{\Theta}\right)\right\}$.

In abstraction, one can consider an isomorphic group where only the commutation rules are considered to be known regarding the operators $T_{i}:\left[T_{1}, T_{2}\right]=i T_{3}$, etc.

In elementary particle physics the $T_{i}$ can be interpreted e.g. as the isospin operators. Elementary particles can be classified in isospin-multiplets, these are the irreducible representations of $\mathrm{SU}(2)$. The classification is:

1. The isospin-singlet $\equiv$ the identical representation: $\mathrm{e}^{-i \vec{T} \cdot \vec{\Theta}}=1 \Rightarrow T_{i}=0$
2. The isospin-doublet $\equiv$ the faithful representation of $\mathrm{SU}(2)$ on $2 \times 2$ matrices.

The group $\mathrm{SU}(3)$ has 8 free parameters. (The group $\mathrm{SU}(N)$ has $N^{2}-1$ free parameters). The Hermitian, traceless operators are $3 \mathrm{SU}(2)$-subgroups in the $\vec{e}_{1} \vec{e}_{2}, \vec{e}_{1} \vec{e}_{3}$ and the $\vec{e}_{2} \vec{e}_{3}$ plane. This gives 9 matrices, which are not all 9 linear independent. By taking a linear combination one gets 8 matrices.
In the Lagrange density for the colour force one has to substitute $\frac{\partial}{\partial x} \rightarrow \frac{D}{D x}:=\frac{\partial}{\partial x}-\sum_{i=1}^{8} T_{i} A_{x}^{i}$
The terms of 3 rd and 4th power in $A$ show that the colour field interacts with itself.

## Chapter 14

## Nuclear physics

### 14.1 Nuclear forces

The mass of a nucleus is given by:

$$
M_{\mathrm{nucl}}=Z m_{\mathrm{p}}+N m_{\mathrm{n}}-E_{\mathrm{bind}} / c^{2}
$$

The binding energy per nucleon is given in the figure at the right. The top is at ${ }_{26}^{56} \mathrm{Fe}$, the most stable nucleus. With the constants

$$
\begin{array}{rlr}
a_{1} & = & 15.760 \\
\mathrm{MeV} \\
a_{2} & = & 17.810 \\
a_{3} & = & \mathrm{MeV} \\
a_{4} & =23.711 & \mathrm{MeV} \\
a_{5} & =34.000 & \mathrm{MeV} \\
\mathrm{MeV}
\end{array}
$$


and $A=Z+N$, in the droplet or collective model of the nucleus the binding energy $E_{\mathrm{bind}}$ is given by:

$$
\frac{E_{\mathrm{bind}}}{c^{2}}=a_{1} A-a_{2} A^{2 / 3}-a_{3} \frac{Z(Z-1)}{A^{1 / 3}}-a_{4} \frac{(N-Z)^{2}}{A}+\epsilon a_{5} A^{-3 / 4}
$$

These terms arise from:

1. $a_{1}$ : Binding energy of the strong nuclear force, approximately $\sim A$.
2. $a_{2}$ : Surface correction: the nucleons near the surface are less bound.
3. $a_{3}$ : Coulomb repulsion between the protons.
4. $a_{4}$ : Asymmetry term: a surplus of protons or neutrons has a lower binding energy.
5. $a_{5}$ : Pair off effect: nuclei with an even number of protons or neutrons are more stable because groups of two protons or neutrons have a lower energy. The following holds:
$\underline{Z}$ even, $N$ even: $\epsilon=+1, \underline{Z}$ odd, $N$ odd: $\epsilon=-1$. $\overline{Z \text { even, } N \text { odd: }} \epsilon=0, \underline{Z \text { odd, } N \text { even: } \epsilon}=0$.

The Yukawa potential can be derived if the nuclear force can to first approximation, be considered as an exchange of virtual pions:

$$
U(r)=-\frac{W_{0} r_{0}}{r} \exp \left(-\frac{r}{r_{0}}\right)
$$

With $\Delta E \cdot \Delta t \approx \hbar, E_{\gamma}=m_{0} c^{2}$ and $r_{0}=c \Delta t$ follows: $r_{0}=\hbar / m_{0} c$.
In the shell model of the nucleus one assumes that a nucleon moves in an average field of other nucleons. Further, there is a contribution of the spin-orbit coupling $\sim \vec{L} \cdot \vec{S}: \Delta V_{l s}=\frac{1}{2}(2 l+1) \hbar \omega$. So each level $(n, l)$ is split in two, with $j=l \pm \frac{1}{2}$, where the state with $j=l+\frac{1}{2}$ has the lowest energy. This is just the opposite for electrons, which is an indication that the $L-S$ interaction is not electromagnetical. The energy of a 3-dimensional harmonic oscillator is $E=\left(N+\frac{3}{2}\right) \hbar \omega . N=n_{x}+n_{y}+n_{z}=2(n-1)+l$ where $n \geq 1$ is the main oscillator number. Because $-l \leq m \leq l$ and $m_{s}= \pm \frac{1}{2} \hbar$ there are $2(2 l+1)$
substates which exist independently for protons and neutrons. This gives rise to the so called magical numbers: nuclei where each state in the outermost level are filled are particulary stable. This is the case if $N$ or $Z$ $\in\{2,8,20,28,50,82,126\}$.

### 14.2 The shape of the nucleus

A nucleus is to first approximation spherical with a radius of $R=R_{0} A^{1 / 3}$. Here, $R_{0} \approx 1.4 \cdot 10^{-15} \mathrm{~m}$, constant for all nuclei. If the nuclear radius is measured including the charge distribution one obtains $R_{0} \approx 1.2 \cdot 10^{-15}$ m . The shape of oscillating nuclei can be described by spherical harmonics:

$$
R=R_{0}\left[1+\sum_{l m} a_{l m} Y_{l}^{m}(\theta, \varphi)\right]
$$

$l=0$ gives rise to monopole vibrations, density vibrations, which can be applied to the theory of neutron stars. $l=1$ gives dipole vibrations, $l=2$ quadrupole, with $a_{2,0}=\beta \cos \gamma$ and $a_{2, \pm 2}=\frac{1}{2} \sqrt{2} \beta \sin \gamma$ where $\beta$ is the deformation factor and $\gamma$ the shape parameter. The multipole moment is given by $\mu_{l}=\operatorname{Zer}^{l} Y_{l}^{m}(\theta, \varphi)$. The parity of the electric moment is $\Pi_{E}=(-1)^{l}$, of the magnetic moment $\Pi_{M}=(-1)^{l+1}$.
There are 2 contributions to the magnetic moment: $\vec{M}_{L}=\frac{e}{2 m_{\mathrm{p}}} \vec{L}$ and $\vec{M}_{S}=g_{S} \frac{e}{2 m_{\mathrm{p}}} \vec{S}$.
where $g_{S}$ is the spin-gyromagnetic ratio. For protons holds $g_{S}=5.5855$ and for neutrons $g_{S}=-3.8263$. The $z$-components of the magnetic moment are given by $M_{L, z}=\mu_{\mathrm{N}} m_{l}$ and $M_{S, z}=g_{S} \mu_{\mathrm{N}} m_{S}$. The resulting magnetic moment is related to the nuclear spin $I$ according to $\vec{M}=g_{I}\left(e / 2 m_{\mathrm{p}}\right) \vec{I}$. The $z$-component is then $M_{z}=\mu_{\mathrm{N}} g_{I} m_{I}$.

### 14.3 Radioactive decay

The number of nuclei decaying is proportional to the number of nuclei: $\dot{N}=-\lambda N$. This gives for the number of nuclei $N: N(t)=N_{0} \exp (-\lambda t)$. The half life time follows from $\tau_{\frac{1}{2}} \lambda=\ln (2)$. The average life time of a nucleus is $\tau=1 / \lambda$. The probability that $N$ nuclei decay within a time interval is given by a Poisson distribution:

$$
P(N) d t=N_{0} \frac{\lambda^{N} \mathrm{e}^{-\lambda}}{N!} d t
$$

If a nucleus can decay into more final states then holds: $\lambda=\sum \lambda_{i}$. So the fraction decaying into state $i$ is $\lambda_{i} / \sum \lambda_{i}$. There are 5 types of natural radioactive decay:

1. $\alpha$-decay: the nucleus emits a $\mathrm{He}^{2+}$ nucleus. Because nucleons tend to order themselves in groups of $2 \mathrm{p}+2 \mathrm{n}$ this can be considered as a tunneling of a $\mathrm{He}^{2+}$ nucleus through a potential barrier. The tunnel probability $P$ is

$$
P=\frac{\text { incoming amplitude }}{\text { outgoing amplitude }}=\mathrm{e}^{-2 G} \text { with } G=\frac{1}{\hbar} \sqrt{2 m \int[V(r)-E] d r}
$$

$G$ is called the Gamow factor.
2. $\beta$-decay. Here a proton changes into a neutron or vice versa:
$\mathrm{p}^{+} \rightarrow \mathrm{n}^{0}+\mathrm{W}^{+} \rightarrow \mathrm{n}^{0}+\mathrm{e}^{+}+\nu_{\mathrm{e}}$, and $\mathrm{n}^{0} \rightarrow \mathrm{p}^{+}+\mathrm{W}^{-} \rightarrow \mathrm{p}^{+}+\mathrm{e}^{-}+\bar{\nu}_{\mathrm{e}}$.
3. Electron capture: here, a proton in the nucleus captures an electron (usually from the K-shell).
4. Spontaneous fission: a nucleus breaks apart.
5. $\gamma$-decay: here the nucleus emits a high-energetic photon. The decay constant is given by

$$
\lambda=\frac{P(l)}{\hbar \omega} \sim \frac{E_{\gamma}}{(\hbar c)^{2}}\left(\frac{E_{\gamma} R}{\hbar c}\right)^{2 l} \sim 10^{-4 l}
$$

where $l$ is the quantum number for the angular momentum and $P$ the radiated power. Usually the decay constant of electric multipole moments is larger than the one of magnetic multipole moments. The energy of the photon is $E_{\gamma}=E_{i}-E_{f}-T_{\mathrm{R}}$, with $T_{\mathrm{R}}=E_{\gamma}^{2} / 2 m c^{2}$ the recoil energy, which can usually be neglected. The parity of the emitted radiation is $\Pi^{l}=\Pi^{i} \cdot \Pi^{f}$. With $I$ the quantum number of angular momentum of the nucleus, $L=\hbar \sqrt{I(I+1)}$, holds the following selection rule: $\left|\vec{I}_{i}-\vec{I}_{f}\right| \leq \Delta l \leq\left|\vec{I}_{i}+\vec{I}_{f}\right|$.

### 14.4 Scattering and nuclear reactions

### 14.4.1 Kinetic model

If a beam with intensity $I$ hits a target with density $n$ and length $x$ (Rutherford scattering) the number of scatterings $R$ per unit of time is equal to $R=\operatorname{Inx} \sigma$. From this follows that the intensity of the beam decreases as $-d I=I n \sigma d x$. This results in $I=I_{0} \mathrm{e}^{-n \sigma x}=I_{0} \mathrm{e}^{-\mu x}$.
Because $d R=R(\theta, \varphi) d \Omega / 4 \pi=I n x d \sigma$ it follows: $\frac{d \sigma}{d \Omega}=\frac{R(\theta, \varphi)}{4 \pi n x I}$
If $N$ particles are scattered in a material with density $n$ then holds: $\frac{\Delta N}{N}=n \frac{d \sigma}{d \Omega} \Delta \Omega \Delta x$
For Coulomb collisions holds: $\left.\frac{d \sigma}{d \Omega}\right|_{\mathrm{C}}=\frac{Z_{1} Z_{2} e^{2}}{8 \pi \varepsilon_{0} \mu v_{0}^{2}} \frac{1}{\sin ^{4}\left(\frac{1}{2} \theta\right)}$

### 14.4.2 Quantum mechanical model for n-p scattering

The initial state is a beam of neutrons moving along the $z$-axis with wavefunction $\psi_{\text {init }}=\mathrm{e}^{i k z}$ and current density $J_{\text {init }}=v\left|\psi_{\text {init }}\right|^{2}=v$. At large distances from the scattering point they have approximately a spherical wavefunction $\psi_{\text {scat }}=f(\theta) \mathrm{e}^{i k r} / r$ where $f(\theta)$ is the scattering amplitude. The total wavefunction is then given by

$$
\psi=\psi_{\text {in }}+\psi_{\text {scat }}=\mathrm{e}^{i k z}+f(\theta) \frac{\mathrm{e}^{i k r}}{r}
$$

The particle flux of the scattered particles is $v\left|\psi_{\text {scat }}\right|^{2}=v|f(\theta)|^{2} d \Omega$. From this it follows that $\sigma(\theta)=|f(\theta)|^{2}$. The wavefunction of the incoming particles can be expressed as a sum of angular momentum wavefunctions:

$$
\psi_{\text {init }}=\mathrm{e}^{i k z}=\sum_{l} \psi_{l}
$$

The impact parameter is related to the angular momentum with $L=b p=b \hbar k$, so $b k \approx l$. At very low energy only particles with $l=0$ are scattered, so

$$
\psi=\psi_{0}^{\prime}+\sum_{l>0} \psi_{l} \text { and } \psi_{0}=\frac{\sin (k r)}{k r}
$$

If the potential is approximately rectangular holds: $\psi_{0}^{\prime}=C \frac{\sin \left(k r+\delta_{0}\right)}{k r}$
The cross section is then $\sigma(\theta)=\frac{\sin ^{2}\left(\delta_{0}\right)}{k^{2}}$ so $\sigma=\int \sigma(\theta) d \Omega=\frac{4 \pi \sin ^{2}\left(\delta_{0}\right)}{k^{2}}$
At very low energies holds: $\sin ^{2}\left(\delta_{0}\right)=\frac{\hbar^{2} k^{2} / 2 m}{W_{0}+W}$
with $W_{0}$ the depth of the potential well. At higher energies holds: $\sigma=\frac{4 \pi}{k^{2}} \sum_{l} \sin ^{2}\left(\delta_{l}\right)$

### 14.4.3 Conservation of energy and momentum in nuclear reactions

If a particle $P_{1}$ collides with a particle $P_{2}$ which is in rest w.r.t. the laboratory system and other particles are created, so

$$
P_{1}+P_{2} \rightarrow \sum_{k>2} P_{k}
$$

the total energy $Q$ gained or required is given by $Q=\left(m_{1}+m_{2}-\sum_{k>2} m_{k}\right) c^{2}$.
The minimal required kinetic energy $T$ of $P_{1}$ in the laboratory system to initialize the reaction is

$$
T=-Q \frac{m_{1}+m_{2}+\sum m_{k}}{2 m_{2}}
$$

If $Q<0$ there is a threshold energy.

### 14.5 Radiation dosimetry

Radiometric quantities determine the strength of the radiation source(s). Dosimetric quantities are related to the energy transfer from radiation to matter. Parameters describing a relation between those are called interaction parameters. The intensity of a beam of particles in matter decreases according to $I(s)=I_{0} \exp (-\mu s)$. The deceleration of a heavy particle is described by the Bethe-Bloch equation:

$$
\frac{d E}{d s} \sim \frac{q^{2}}{v^{2}}
$$

The fluention is given by $\Phi=d N / d A$. The flux is given by $\phi=d \Phi / d t$. The energy loss is defined by $\Psi=$ $d W / d A$, and the energy flux density $\psi=d \Psi / d t$. The absorption coefficient is given by $\mu=(d N / N) / d x$. The mass absorption coefficient is given by $\mu / \varrho$.
The radiation dose $X$ is the amount of charge produced by the radiation per unit of mass, with unit $\mathrm{C} / \mathrm{kg}$. An old unit is the Röntgen: $1 \mathrm{Ro}=2.58 \cdot 10^{-4} \mathrm{C} / \mathrm{kg}$. With the energy-absorption coefficient $\mu_{E}$ follows:

$$
X=\frac{d Q}{d m}=\frac{e \mu_{E}}{W \varrho} \Psi
$$

where $W$ is the energy required to disjoin an elementary charge.
The absorbed dose $D$ is given by $D=d E_{\text {abs }} / d m$, with unit $G y=\mathrm{J} / \mathrm{kg}$. An old unit is the rad: $1 \mathrm{rad}=0.01 \mathrm{~Gy}$. The dose tempo is defined as $\dot{D}$. It can be derived that

$$
D=\frac{\mu_{E}}{\varrho} \Psi
$$

The Kerma $K$ is the amount of kinetic energy of secundary produced particles which is produced per mass unit of the radiated object.

The equivalent dose $H$ is a weight average of the absorbed dose per type of radiation, where for each type radiation the effects on biological material is used for the weight factor. These weight factors are called the quality factors. Their unit is $\operatorname{Sv} . H=Q D$. If the absorption is not equally distributed also weight factors $w$ per organ need to be used: $H=\sum w_{k} H_{k}$. For some types of radiation holds:

| Radiation type | $\boldsymbol{Q}$ |
| :--- | ---: |
| Röntgen, gamma radiation | 1 |
| $\beta$, electrons, mesons | 1 |
| Thermic neutrons | 3 to 5 |
| Fast neutrons | 10 to 20 |
| protons | 10 |
| $\alpha$, fission products | 20 |

## Chapter 15

## Quantum fi eld theory \& Particle physics

### 15.1 Creation and annihilation operators

A state with more particles can be described by a collection occupation numbers $\left|n_{1} n_{2} n_{3} \cdots\right\rangle$. Hence the vacuum state is given by $|000 \cdots\rangle$. This is a complete description because the particles are indistinguishable. The states are orthonormal:

$$
\left\langle n_{1} n_{2} n_{3} \cdots \mid n_{1}^{\prime} n_{2}^{\prime} n_{3}^{\prime} \cdots\right\rangle=\prod_{i=1}^{\infty} \delta_{n_{i} n_{i}^{\prime}}
$$

The time-dependent state vector is given by

$$
\Psi(t)=\sum_{n_{1} n_{2} \cdots} c_{n_{1} n_{2} \cdots}(t)\left|n_{1} n_{2} \cdots\right\rangle
$$

The coefficients $c$ can be interpreted as follows: $\left|c_{n_{1} n_{2} \ldots}\right|^{2}$ is the probability to find $n_{1}$ particles with momentum $\vec{k}_{1}, n_{2}$ particles with momentum $\vec{k}_{2}$, etc., and $\langle\Psi(t) \mid \Psi(t)\rangle=\sum\left|c_{n_{i}}(t)\right|^{2}=1$. The expansion of the states in time is described by the Schrödinger equation

$$
i \frac{d}{d t}|\Psi(t)\rangle=H|\Psi(t)\rangle
$$

where $H=H_{0}+H_{\mathrm{int}}$. $H_{0}$ is the Hamiltonian for free particles and keeps $\left|c_{n_{i}}(t)\right|^{2}$ constant, $H_{\text {int }}$ is the interaction Hamiltonian and can increase or decrease a $c^{2}$ at the cost of others.
All operators which can change occupation numbers can be expanded in the $a$ and $a^{\dagger}$ operators. $a$ is the annihilation operator and $a^{\dagger}$ the creation operator, and:

$$
\begin{aligned}
a\left(\vec{k}_{i}\right)\left|n_{1} n_{2} \cdots n_{i} \cdots\right\rangle & =\sqrt{n_{i}}\left|n_{1} n_{2} \cdots n_{i}-1 \cdots\right\rangle \\
a^{\dagger}\left(\vec{k}_{i}\right)\left|n_{1} n_{2} \cdots n_{i} \cdots\right\rangle & =\sqrt{n_{i}+1}\left|n_{1} n_{2} \cdots n_{i}+1 \cdots\right\rangle
\end{aligned}
$$

Because the states are normalized holds $a|0\rangle=0$ and $a\left(\vec{k}_{i}\right) a^{\dagger}\left(\vec{k}_{i}\right)\left|n_{i}\right\rangle=n_{i}\left|n_{i}\right\rangle$. So $a a^{\dagger}$ is an occupation number operator. The following commutation rules can be derived:

$$
\left[a\left(\vec{k}_{i}\right), a\left(\vec{k}_{j}\right)\right]=0, \quad\left[a^{\dagger}\left(\vec{k}_{i}\right), a^{\dagger}\left(\vec{k}_{j}\right)\right]=0 \quad, \quad\left[a\left(\vec{k}_{i}\right), a^{\dagger}\left(\vec{k}_{j}\right)\right]=\delta_{i j}
$$

Hence for free spin-0 particles holds: $H_{0}=\sum_{i} a^{\dagger}\left(\vec{k}_{i}\right) a\left(\vec{k}_{i}\right) \hbar \omega_{k_{i}}$

### 15.2 Classical and quantum fields

Starting with a real field $\Phi^{\alpha}(x)$ (complex fields can be split in a real and an imaginary part), the Lagrange density $\mathcal{L}$ is a function of the position $x=(\vec{x}, i c t)$ through the fields: $\mathcal{L}=\mathcal{L}\left(\Phi^{\alpha}(x), \partial_{\nu} \Phi^{\alpha}(x)\right)$. The Lagrangian is given by $L=\int \mathcal{L}(x) d^{3} x$. Using the variational principle $\delta I(\Omega)=0$ and with the action-integral $I(\Omega)=\int \mathcal{L}\left(\Phi^{\alpha}, \partial_{\nu} \Phi^{\alpha}\right) d^{4} x$ the field equation can be derived:

$$
\frac{\partial \mathcal{L}}{\partial \Phi^{\alpha}}-\frac{\partial}{\partial x_{\nu}} \frac{\partial \mathcal{L}}{\partial\left(\partial_{\nu} \Phi^{\alpha}\right)}=0
$$

The conjugated field is, analogous to momentum in classical mechanics, defined as:

$$
\Pi^{\alpha}(x)=\frac{\partial \mathcal{L}}{\partial \dot{\Phi}^{\alpha}}
$$

With this, the Hamilton density becomes $\mathcal{H}(x)=\Pi^{\alpha} \dot{\Phi}^{\alpha}-\mathcal{L}(x)$.
Quantization of a classical field is analogous to quantization in point mass mechanics: the field functions are considered as operators obeying certain commutation rules:

$$
\left[\Phi^{\alpha}(\vec{x}), \Phi^{\beta}\left(\vec{x}^{\prime}\right)\right]=0, \quad\left[\Pi^{\alpha}(\vec{x}), \Pi^{\beta}\left(\vec{x}^{\prime}\right)\right]=0 \quad, \quad\left[\Phi^{\alpha}(\vec{x}), \Pi^{\beta}\left(\vec{x}^{\prime}\right)\right]=i \delta_{\alpha \beta}\left(\vec{x}-\vec{x}^{\prime}\right)
$$

### 15.3 The interaction picture

Some equivalent formulations of quantum mechanics are possible:

1. Schrödinger picture: time-dependent states, time-independent operators.
2. Heisenberg picture: time-independent states, time-dependent operators.
3. Interaction picture: time-dependent states, time-dependent operators.

The interaction picture can be obtained from the Schrödinger picture by an unitary transformation:

$$
|\Phi(t)\rangle=\mathrm{e}^{i H_{0}^{\mathrm{S}}}\left|\Phi^{\mathrm{S}}(t)\right\rangle \text { and } O(t)=\mathrm{e}^{i H_{0}^{\mathrm{S}}} O^{\mathrm{S}} \mathrm{e}^{-i H_{0}^{\mathrm{S}}}
$$

The index ${ }^{\mathrm{S}}$ denotes the Schrödinger picture. From this follows:

$$
i \frac{d}{d t}|\Phi(t)\rangle=H_{\mathrm{int}}(t)|\Phi(t)\rangle \text { and } i \frac{d}{d t} O(t)=\left[O(t), H_{0}\right]
$$

### 15.4 Real scalar field in the interaction picture

It is easy to find that, with $M:=m_{0}^{2} c^{2} / \hbar^{2}$, holds:

$$
\frac{\partial}{\partial t} \Phi(x)=\Pi(x) \text { and } \frac{\partial}{\partial t} \Pi(x)=\left(\nabla^{2}-M^{2}\right) \Phi(x)
$$

From this follows that $\Phi$ obeys the Klein-Gordon equation $\left(\square-M^{2}\right) \Phi=0$. With the definition $k_{0}^{2}=$ $\vec{k}^{2}+M^{2}:=\omega_{k}^{2}$ and the notation $\vec{k} \cdot \vec{x}-i k_{0} t:=k x$ the general solution of this equation is:
$\Phi(x)=\frac{1}{\sqrt{V}} \sum_{\vec{k}} \frac{1}{\sqrt{2 \omega_{k}}}\left(a(\vec{k}) \mathrm{e}^{i k x}+a^{\dagger}(\vec{k}) \mathrm{e}^{-i k x}\right), \Pi(x)=\frac{i}{\sqrt{V}} \sum_{\vec{k}} \sqrt{\frac{1}{2} \omega_{k}}\left(-a(\vec{k}) \mathrm{e}^{i k x}+a^{\dagger}(\vec{k}) \mathrm{e}^{-i k x}\right)$
The field operators contain a volume $V$, which is used as normalization factor. Usually one can take the limit $V \rightarrow \infty$.
In general it holds that the term with $\mathrm{e}^{-i k x}$, the positive frequency part, is the creation part, and the negative frequency part is the annihilation part.
the coefficients have to be each others hermitian conjugate because $\Phi$ is hermitian. Because $\Phi$ has only one component this can be interpreted as a field describing a particle with spin zero. From this follows that the commutation rules are given by $\left[\Phi(x), \Phi\left(x^{\prime}\right)\right]=i \Delta\left(x-x^{\prime}\right)$ with

$$
\Delta(y)=\frac{1}{(2 \pi)^{3}} \int \frac{\sin (k y)}{\omega_{k}} d^{3} k
$$

$\Delta(y)$ is an odd function which is invariant for proper Lorentz transformations (no mirroring). This is consistent with the previously found result $\left[\Phi\left(\vec{x}, t, \Phi\left(\vec{x}^{\prime}, t\right)\right]=0\right.$. In general holds that $\Delta(y)=0$ outside the light cone. So the equations obey the locality postulate.
The Lagrange density is given by: $\mathcal{L}\left(\Phi, \partial_{\nu} \Phi\right)=-\frac{1}{2}\left(\partial_{\nu} \Phi \partial_{\nu} \Phi+m^{2} \Phi^{2}\right)$. The energy operator is given by:

$$
H=\int \mathcal{H}(x) d^{3} x=\sum_{\vec{k}} \hbar \omega_{k} a^{\dagger}(\vec{k}) a(\vec{k})
$$

### 15.5 Charged spin-0 particles, conservation of charge

The Lagrange density of charged spin-0 particles is given by: $\mathcal{L}=-\left(\partial_{\nu} \Phi \partial_{\nu} \Phi^{*}+M^{2} \Phi \Phi^{*}\right)$.
Noether's theorem connects a continuous symmetry of $\mathcal{L}$ and an additive conservation law. Suppose that $\mathcal{L}\left(\left(\Phi^{\alpha}\right)^{\prime}, \partial_{\nu}\left(\Phi^{\alpha}\right)^{\prime}\right)=\mathcal{L}\left(\Phi^{\alpha}, \partial_{\nu} \Phi^{\alpha}\right)$ and there exists a continuous transformation between $\Phi^{\alpha}$ and $\Phi^{\alpha \prime}$ such as $\Phi^{\alpha \prime}=\Phi^{\alpha}+\epsilon f^{\alpha}(\Phi)$. Then holds

$$
\frac{\partial}{\partial x_{\nu}}\left(\frac{\partial \mathcal{L}}{\partial\left(\partial_{\nu} \Phi^{\alpha}\right)} f^{\alpha}\right)=0
$$

This is a continuity equation $\Rightarrow$ conservation law. Which quantity is conserved depends on the symmetry. The above Lagrange density is invariant for a change in phase $\Phi \rightarrow \Phi \mathrm{e}^{i \theta}$ : a global gauge transformation. The conserved quantity is the current density $J_{\mu}(x)=-i e\left(\Phi \partial_{\mu} \Phi^{*}-\Phi^{*} \partial_{\mu} \Phi\right)$. Because this quantity is 0 for real fields a complex field is needed to describe charged particles. When this field is quantized the field operators are given by

$$
\Phi(x)=\frac{1}{\sqrt{V}} \sum_{\vec{k}} \frac{1}{\sqrt{2 \omega_{k}}}\left(a(\vec{k}) \mathrm{e}^{i k x}+b^{\dagger}(\vec{k}) \mathrm{e}^{-i k x}\right), \Phi^{\dagger}(x)=\frac{1}{\sqrt{V}} \sum_{\vec{k}} \frac{1}{\sqrt{2 \omega_{k}}}\left(a^{\dagger}(\vec{k}) \mathrm{e}^{i k x}+b(\vec{k}) \mathrm{e}^{-i k x}\right)
$$

Hence the energy operator is given by:

$$
H=\sum_{\vec{k}} \hbar \omega_{k}\left(a^{\dagger}(\vec{k}) a(\vec{k})+b^{\dagger}(\vec{k}) b(\vec{k})\right)
$$

and the charge operator is given by:

$$
Q(t)=-i \int J_{4}(x) d^{3} x \Rightarrow Q=\sum_{\vec{k}} e\left(a^{\dagger}(\vec{k}) a(\vec{k})-b^{\dagger}(\vec{k}) b(\vec{k})\right)
$$

From this follows that $a^{\dagger} a:=N_{+}(\vec{k})$ is an occupation number operator for particles with a positive charge and $b^{\dagger} b:=N_{-}(\vec{k})$ is an occupation number operator for particles with a negative charge.

### 15.6 Field functions for spin- $\frac{1}{2}$ particles

Spin is defined by the behaviour of the solutions $\psi$ of the Dirac equation. A scalar field $\Phi$ has the property that, if it obeys the Klein-Gordon equation, the rotated field $\tilde{\Phi}(x):=\Phi\left(\Lambda^{-1} x\right)$ also obeys it. $\Lambda$ denotes 4-dimensional rotations: the proper Lorentz transformations. These can be written as:

$$
\tilde{\Phi}(x)=\Phi(x) \mathrm{e}^{-i \vec{n} \cdot \vec{L}} \text { with } L_{\mu \nu}=-i \hbar\left(x_{\mu} \frac{\partial}{\partial x_{\nu}}-x_{\nu} \frac{\partial}{\partial x_{\mu}}\right)
$$

For $\mu \leq 3, \nu \leq 3$ these are rotations, for $\nu=4, \mu \neq 4$ these are Lorentz transformations.
A rotated field $\tilde{\psi}$ obeys the Dirac equation if the following condition holds: $\tilde{\psi}(x)=D(\Lambda) \psi\left(\Lambda^{-1} x\right)$. This results in the condition $D^{-1} \gamma_{\lambda} D=\Lambda_{\lambda \mu} \gamma_{\mu}$. One finds: $D=\mathrm{e}^{i \vec{n} \cdot \vec{S}}$ with $S_{\mu \nu}=-i \frac{1}{2} \hbar \gamma_{\mu} \gamma_{\nu}$. Hence:

$$
\tilde{\psi}(x)=\mathrm{e}^{-i(S+L)} \psi(x)=\mathrm{e}^{-i J} \psi(x)
$$

Then the solutions of the Dirac equation are given by:

$$
\psi(x)=u_{ \pm}^{r}(\vec{p}) \mathrm{e}^{-i(\vec{p} \cdot \vec{x} \pm E t)}
$$

Here, $r$ is an indication for the direction of the spin, and $\pm$ is the sign of the energy. With the notation $v^{r}(\vec{p})=u_{-}^{r}(-\vec{p})$ and $u^{r}(\vec{p})=u_{+}^{r}(\vec{p})$ one can write for the dot products of these spinors:

$$
u_{+}^{r}(\vec{p}) u_{+}^{r^{\prime}}(\vec{p})=\frac{E}{M} \delta_{r r^{\prime}}, u_{-}^{r}(\vec{p}) u_{-}^{r^{\prime}}(\vec{p})=\frac{E}{M} \delta_{r r^{\prime}}, u_{+}^{r}(\vec{p}) u_{-}^{r^{\prime}}(\vec{p})=0
$$

Because of the factor $E / M$ this is not relativistic invariant. A Lorentz-invariant dot product is defined by $\bar{a} b:=a^{\dagger} \gamma_{4} b$, where $\bar{a}:=a^{\dagger} \gamma_{4}$ is a row spinor. From this follows:

$$
\overline{u^{r}(\vec{p})} u^{r^{\prime}}(\vec{p})=\delta_{r r^{\prime}}, \overline{v^{r}(\vec{p})} v^{r^{\prime}}(\vec{p})=-\delta_{r r^{\prime}}, \overline{u^{r}(\vec{p})} v^{r^{\prime}}(\vec{p})=0
$$

Combinations of the type $a \bar{a}$ give a $4 \times 4$ matrix:

$$
\sum_{r=1}^{2} u^{r}(\vec{p}) \overline{u^{r}(\vec{p})}=\frac{-i \gamma_{\lambda} p_{\lambda}+M}{2 M}, \quad \sum_{r=1}^{2} v^{r}(\vec{p}) \overline{v^{r}(\vec{p})}=\frac{-i \gamma_{\lambda} p_{\lambda}-M}{2 M}
$$

The Lagrange density which results in the Dirac equation and having the correct energy normalization is:

$$
\mathcal{L}(x)=-\overline{\psi(x)}\left(\gamma_{\mu} \frac{\partial}{\partial x_{\mu}}+M\right) \psi(x)
$$

and the current density is $J_{\mu}(x)=-i e \bar{\psi} \gamma_{\mu} \psi$.

### 15.7 Quantization of spin- $\frac{1}{2}$ fields

The general solution for the fieldoperators is in this case:

$$
\psi(x)=\sqrt{\frac{M}{V}} \sum_{\vec{p}} \frac{1}{\sqrt{E}} \sum_{r}\left(c_{r}(\vec{p}) u^{r}(\vec{p}) \mathrm{e}^{i p x}+d_{r}^{\dagger}(\vec{p}) v^{r}(\vec{p}) \mathrm{e}^{-i p x}\right)
$$

and

$$
\overline{\psi(x)}=\sqrt{\frac{M}{V}} \sum_{\vec{p}} \frac{1}{\sqrt{E}} \sum_{r}\left(c_{r}^{\dagger}(\vec{p}) \overline{u^{r}}(\vec{p}) \mathrm{e}^{-i p x}+d_{r}(\vec{p}) \overline{v^{r}}(\vec{p}) \mathrm{e}^{i p x}\right)
$$

Here, $c^{\dagger}$ and $c$ are the creation respectively annihilation operators for an electron and $d^{\dagger}$ and $d$ the creation respectively annihilation operators for a positron. The energy operator is given by

$$
H=\sum_{\vec{p}} E_{\vec{p}} \sum_{r=1}^{2}\left(c_{r}^{\dagger}(\vec{p}) c_{r}(\vec{p})-d_{r}(\vec{p}) d_{r}^{\dagger}(\vec{p})\right)
$$

To prevent that the energy of positrons is negative the operators must obey anti commutation rules in stead of commutation rules:

$$
\left[c_{r}(\vec{p}), c_{r^{\prime}}^{\dagger}(\vec{p})\right]_{+}=\left[d_{r}(\vec{p}), d_{r^{\prime}}^{\dagger}(\vec{p})\right]_{+}=\delta_{r r^{\prime}} \delta_{p p^{\prime}}, \text { all other anti commutators are } 0
$$

The field operators obey

$$
\begin{gathered}
{\left[\psi_{\alpha}(x), \psi_{\beta}\left(x^{\prime}\right)\right]=0,\left[\overline{\psi_{\alpha}(x)}, \overline{\psi_{\beta}\left(x^{\prime}\right)}\right]=0, \quad\left[\psi_{\alpha}(x), \overline{\psi_{\beta}\left(x^{\prime}\right)}\right]_{+}=-i S_{\alpha \beta}\left(x-x^{\prime}\right)} \\
\text { with } S(x)=\left(\gamma_{\lambda} \frac{\partial}{\partial x_{\lambda}}-M\right) \Delta(x)
\end{gathered}
$$

The anti commutation rules give besides the positive-definite energy also the Pauli exclusion principle and the Fermi-Dirac statistics: because $c_{r}^{\dagger}(\vec{p}) c_{r}^{\dagger}(\vec{p})=-c_{r}^{\dagger}(\vec{p}) c_{r}^{\dagger}(\vec{p})$ holds: $\left\{c_{r}^{\dagger}(p)\right\}^{2}=0$. It appears to be impossible to create two electrons with the same momentum and spin. This is the exclusion principle. Another way to see this is the fact that $\left\{N_{r}^{+}(\vec{p})\right\}^{2}=N_{r}^{+}(\vec{p})$ : the occupation operators have only eigenvalues 0 and 1 .
To avoid infinite vacuum contributions to the energy and charge the normal product is introduced. The expression for the current density now becomes $J_{\mu}=-i e N\left(\bar{\psi} \gamma_{\mu} \psi\right)$. This product is obtained by:

- Expand all fields into creation and annihilation operators,
- Keep all terms which have no annihilation operators, or in which they are on the right of the creation operators,
- In all other terms interchange the factors so that the annihilation operators go to the right. By an interchange of two fermion operators add a minus sign, by interchange of two boson operators not. Assume hereby that all commutators are zero.


### 15.8 Quantization of the electromagnetic field

Starting with the Lagrange density $\mathcal{L}=-\frac{1}{2} \frac{\partial A_{\nu}}{\partial x_{\mu}} \frac{\partial A_{\nu}}{\partial x_{\mu}}$
it follows for the field operators $A(x)$ :

$$
A(x)=\frac{1}{\sqrt{V}} \sum_{\vec{k}} \frac{1}{\sqrt{2 \omega_{k}}} \sum_{m=1}^{4}\left(a_{m}(\vec{k}) \epsilon^{m}(\vec{k}) \mathrm{e}^{i k x}+a^{\dagger}(\vec{k}) \epsilon^{m}(\vec{k})^{*} \mathrm{e}^{-i k x}\right)
$$

The operators obey $\left[a_{m}(\vec{k}), a_{m^{\prime}}^{\dagger}(\vec{k})\right]=\delta_{m m^{\prime}} \delta_{k k^{\prime}}$. All other commutators are 0 . $m$ gives the polarization direction of the photon: $m=1,2$ gives transversal polarized, $m=3$ longitudinal polarized and $m=4$ timelike polarized photons. Further holds:

$$
\left[A_{\mu}(x), A_{\nu}\left(x^{\prime}\right)\right]=i \delta_{\mu \nu} D\left(x-x^{\prime}\right) \text { with } D(y)=\left.\Delta(y)\right|_{m=0}
$$

In spite of the fact that $A_{4}=i V$ is imaginary in the classical case, $A_{4}$ is still defined to be hermitian because otherwise the sign of the energy becomes incorrect. By changing the definition of the inner product in configuration space the expectation values for $A_{1,2,3}(x) \in \mathbb{R}$ and for $A_{4}(x)$ become imaginary.
If the potentials satisfy the Lorentz gauge condition $\partial_{\mu} A_{\mu}=0$ the $E$ and $B$ operators derived from these potentials will satisfy the Maxwell equations. However, this gives problems with the commutation rules. It is now demanded that only those states are permitted for which holds

$$
\frac{\partial A_{\mu}^{+}}{\partial x_{\mu}}|\Phi\rangle=0
$$

This results in:

$$
\left\langle\frac{\partial A_{\mu}}{\partial x_{\mu}}\right\rangle=0
$$

From this follows that $\left(a_{3}(\vec{k})-a_{4}(\vec{k})\right)|\Phi\rangle=0$. With a local gauge transformation one obtains $N_{3}(\vec{k})=0$ and $N_{4}(\vec{k})=0$. However, this only applies to free EM-fields: in intermediary states in interactions there can exist longitudinal and timelike photons. These photons are also responsible for the stationary Coulomb potential.

### 15.9 Interacting fields and the S-matrix

The $S$ (scattering)-matrix gives a relation between the initial and final states of an interaction: $|\Phi(\infty)\rangle=$ $S|\Phi(-\infty)\rangle$. If the Schrödinger equation is integrated:

$$
|\Phi(t)\rangle=|\Phi(-\infty)\rangle-i \int_{-\infty}^{t} H_{\mathrm{int}}\left(t_{1}\right)\left|\Phi\left(t_{1}\right)\right\rangle d t_{1}
$$

and perturbation theory is applied one finds that:

$$
S=\sum_{n=0}^{\infty} \frac{(-i)^{n}}{n!} \int \cdots \int T\left\{\mathcal{H}_{\mathrm{int}}\left(x_{1}\right) \cdots \mathcal{H}_{\mathrm{int}}\left(x_{n}\right)\right\} d^{4} x_{1} \cdots d^{4} x_{n} \equiv \sum_{n=0}^{\infty} S^{(n)}
$$

Here, the $T$-operator means a time-ordered product: the terms in such a product must be ordered in increasing time order from the right to the left so that the earliest terms act first. The $S$-matrix is then given by: $S_{i j}=$ $\left\langle\Phi_{i}\right| S\left|\Phi_{j}\right\rangle=\left\langle\Phi_{i} \mid \Phi(\infty)\right\rangle$.
The interaction Hamilton density for the interaction between the electromagnetic and the electron-positron field is: $\mathcal{H}_{\text {int }}(x)=-J_{\mu}(x) A_{\mu}(x)=i e N\left(\bar{\psi} \gamma_{\mu} \psi A_{\mu}\right)$
When this is expanded as: $\mathcal{H}_{\mathrm{int}}=i e N\left(\left(\overline{\psi^{+}}+\overline{\psi^{-}}\right) \gamma_{\mu}\left(\psi^{+}+\psi^{-}\right)\left(A_{\mu}^{+}+A_{\mu}^{-}\right)\right)$
eight terms appear. Each term corresponds with a possible process. The term $i e \overline{\psi^{+}} \gamma_{\mu} \psi^{+} A_{\mu}^{-}$acting on $|\Phi\rangle$ gives transitions where $A_{\mu}^{-}$creates a photon, $\psi^{+}$annihilates an electron and $\overline{\psi^{+}}$annihilates a positron. Only terms with the correct number of particles in the initial and final state contribute to a matrix element $\left\langle\Phi_{i}\right| S\left|\Phi_{j}\right\rangle$. Further the factors in $\mathcal{H}_{\text {int }}$ can create and thereafter annihilate particles: the virtual particles.
The expressions for $S^{(n)}$ contain time-ordered products of normal products. This can be written as a sum of normal products. The appearing operators describe the minimal changes necessary to change the initial state into the final state. The effects of the virtual particles are described by the (anti)commutator functions. Some time-ordened products are:

$$
\begin{aligned}
T\{\Phi(x) \Phi(y)\} & =N\{\Phi(x) \Phi(y)\}+\frac{1}{2} \Delta^{\mathrm{F}}(x-y) \\
T\left\{\psi_{\alpha}(x) \overline{\psi_{\beta}(y)}\right\} & =N\left\{\psi_{\alpha}(x) \overline{\psi_{\beta}(y)}\right\}-\frac{1}{2} S_{\alpha \beta}^{\mathrm{F}}(x-y) \\
T\left\{A_{\mu}(x) A_{\nu}(y)\right\} & =N\left\{A_{\mu}(x) A_{\nu}(y)\right\}+\frac{1}{2} \delta_{\mu \nu} D_{\mu \nu}^{\mathrm{F}}(x-y)
\end{aligned}
$$

Here, $S^{\mathrm{F}}(x)=\left(\gamma_{\mu} \partial_{\mu}-M\right) \Delta^{\mathrm{F}}(x), D^{\mathrm{F}}(x)=\left.\Delta^{\mathrm{F}}(x)\right|_{m=0}$ and

$$
\Delta^{\mathrm{F}}(x)= \begin{cases}\frac{1}{(2 \pi)^{3}} \int \frac{\mathrm{e}^{i k x}}{\omega_{\vec{k}}} d^{3} k & \text { if } x_{0}>0 \\ \frac{1}{(2 \pi)^{3}} \int \frac{\mathrm{e}^{-i k x}}{\omega_{\vec{k}}} d^{3} k & \text { if } x_{0}<0\end{cases}
$$

The term $\frac{1}{2} \Delta^{\mathrm{F}}(x-y)$ is called the contraction of $\Phi(x)$ and $\Phi(y)$, and is the expectation value of the timeordered product in the vacuum state. Wick's theorem gives an expression for the time-ordened product of an arbitrary number of field operators. The graphical representation of these processes are called Feynman diagrams. In the $x$-representation each diagram describes a number of processes. The contraction functions can also be written as:

$$
\Delta^{\mathrm{F}}(x)=\lim _{\epsilon \rightarrow 0} \frac{-2 i}{(2 \pi)^{4}} \int \frac{\mathrm{e}^{i k x}}{k^{2}+m^{2}-i \epsilon} d^{4} k \quad \text { and } \quad S^{\mathrm{F}}(x)=\lim _{\epsilon \rightarrow 0} \frac{-2 i}{(2 \pi)^{4}} \int \mathrm{e}^{i p x} \frac{i \gamma_{\mu} p_{\mu}-M}{p^{2}+M^{2}-i \epsilon} d^{4} p
$$

In the expressions for $S^{(2)}$ this gives rise to terms $\delta\left(p+k-p^{\prime}-k^{\prime}\right)$. This means that energy and momentum is conserved. However, virtual particles do not obey the relation between energy and momentum.

### 15.10 Divergences and renormalization

It turns out that higher orders contribute infinite terms because only the sum $p+k$ of the four-momentum of the virtual particles is fixed. An integration over one of them becomes $\infty$. In the $x$-representation this can be understood because the product of two functions containing $\delta$-like singularities is not well defined. This is solved by discounting all divergent diagrams in a renormalization of $e$ and $M$. It is assumed that an electron, if there would not be an electromagnetical field, would have a mass $M_{0}$ and a charge $e_{0}$ unequal to the observed mass $M$ and charge $e$. In the Hamilton and Lagrange density of the free electron-positron field appears $M_{0}$. So this gives, with $M=M_{0}+\Delta M$ :

$$
\mathcal{L}_{\mathrm{e}-\mathrm{p}}(x)=-\overline{\psi(x)}\left(\gamma_{\mu} \partial_{\mu}+M_{0}\right) \psi(x)=-\overline{\psi(x)}\left(\gamma_{\mu} \partial_{\mu}+M\right) \psi(x)+\Delta M \overline{\psi(x)} \psi(x)
$$

and $\mathcal{H}_{\mathrm{int}}=i e N\left(\bar{\psi} \gamma_{\mu} \psi A_{\mu}\right)-i \Delta e N\left(\bar{\psi} \gamma_{\mu} \psi A_{\mu}\right)$.

### 15.11 Classification of elementary particles

Elementary particles can be categorized as follows:

1. Hadrons: these exist of quarks and can be categorized in:
I. Baryons: these exist of 3 quarks or 3 antiquarks.
II. Mesons: these exist of one quark and one antiquark.
2. Leptons: $\mathrm{e}^{ \pm}, \mu^{ \pm}, \tau^{ \pm}, \nu_{\mathrm{e}}, \nu_{\mu}, \nu_{\tau}, \bar{\nu}_{\mathrm{e}}, \bar{\nu}_{\mu}, \bar{\nu}_{\tau}$.
3. Field quanta: $\gamma, \mathbf{W}^{ \pm}, Z^{0}$, gluons, gravitons (?).

An overview of particles and antiparticles is given in the following table:

| Particle | spin $(\hbar) \mathrm{B}$ | L | T | $\mathrm{T}_{3}$ | S | C | $\mathrm{B}^{*}$ | charge $(e)$ | $m_{0}(\mathrm{MeV})$ | antipart. |
| :---: | ---: | :---: | :---: | :---: | :---: | :---: | ---: | ---: | ---: | :---: |
| u | $1 / 2$ | $1 / 3$ | 0 | $1 / 2$ | $1 / 2$ | 0 | 0 | 0 | $+2 / 3$ | 5 |
| d | $1 / 2$ | $1 / 3$ | 0 | $1 / 2$ | $-1 / 2$ | 0 | 0 | 0 | $-1 / 3$ | $\overline{\mathrm{u}}$ |
| s | $1 / 2$ | $1 / 3$ | 0 | 0 | 0 | -1 | 0 | 0 | $-1 / 3$ | 175 |
| c | $1 / 2$ | $1 / 3$ | 0 | 0 | 0 | 0 | 1 | 0 | $+2 / 3$ | 1350 |
| d |  |  |  |  |  |  |  |  |  |  |
| b | $1 / 2$ | $1 / 3$ | 0 | 0 | 0 | 0 | 0 | -1 | $-1 / 3$ | 4500 |
| s | $1 / 2$ | $1 / 3$ | 0 | 0 | 0 | 0 | 0 | 0 | $+2 / 3$ | 173000 |
| $\overline{\mathrm{c}}$ |  |  |  |  |  |  |  |  |  |  |
| $\mathrm{e}^{-}$ | $1 / 2$ | 0 | 1 | 0 | 0 | 0 | 0 | 0 | -1 | 0.511 |
| $\mu^{-}$ | $1 / 2$ | 0 | 1 | 0 | 0 | 0 | 0 | 0 | -1 | 105.658 |
| $\tau^{-}$ | $1 / 2$ | 0 | 1 | 0 | 0 | 0 | 0 | 0 | -1 | 1777.1 |
| $\nu_{\mathrm{e}}$ | $1 / 2$ | 0 | 1 | 0 | 0 | 0 | 0 | 0 | $\mathrm{e}^{+}$ |  |
| $\nu_{\mu}$ | $1 / 2$ | 0 | 1 | 0 | 0 | 0 | 0 | 0 | $\tau^{+}$ |  |
| $\nu_{\tau}$ | $1 / 2$ | 0 | 1 | 0 | 0 | 0 | 0 | 0 | 0 | $0(?)$ |
| $\gamma$ | 1 | 0 | 0 | 0 | 0 | 0 | 0 | 0 | 0 | $\bar{\nu}_{\mathrm{e}}$ |
| gluon | 1 | 0 | 0 | 0 | 0 | 0 | 0 | 0 | 0 | $\bar{\nu}_{\mu}$ |
| $\mathrm{W}^{+}$ | 1 | 0 | 0 | 0 | 0 | 0 | 0 | 0 | 0 | +1 |
| Z | 1 | 0 | 0 | 0 | 0 | 0 | 0 | 0 | 80 | $\bar{\nu}_{\tau}$ |
| graviton | 2 | 0 | 0 | 0 | 0 | 0 | 0 | 0 | 0 | 0 |

Here B is the baryon number and L the lepton number. It is found that there are three different lepton numbers, one for $\mathrm{e}, \mu$ and $\tau$, which are separately conserved. T is the isospin, with $T_{3}$ the projection of the isospin on the third axis, C the charmness, S the strangeness and $\mathrm{B}^{*}$ the bottomness. The anti particles have quantum numbers with the opposite sign except for the total isospin $T$. The composition of (anti)quarks of the hadrons is given in the following table, together with their mass in MeV in their ground state:

| $\pi^{0}$ | $\frac{1}{2} \sqrt{2}(\mathrm{u} \overline{\mathrm{u}}+\mathrm{d} \overline{\mathrm{d}})$ | 134.9764 | $\mathrm{J} / \Psi$ | c $\bar{C}$ | 3096.8 | $\overline{\Sigma^{+}}$ | $\overline{\mathrm{d}} \overline{\mathrm{d}} \overline{\mathrm{s}}$ | 1197.436 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\pi^{+}$ | u $\bar{d}$ | 139.56995 | $\Upsilon$ | $\mathrm{b} \overline{\mathrm{b}}$ | 9460.37 | $\Xi^{0}$ | us s | 1314.9 |
| $\pi^{-}$ | d $\bar{u}$ | 139.56995 | $\mathrm{p}^{+}$ | u ud | 938.27231 | $\bar{\Xi}^{0}$ | $\overline{\mathrm{u}} \overline{\mathrm{s}} \overline{\mathrm{s}}$ | 1314.9 |
| $\mathrm{K}^{0}$ | s $\bar{d}$ | 497.672 | $\mathrm{p}^{-}$ | $\bar{u} \bar{u} \bar{d}$ | 938.27231 | $\Xi^{-}$ | d s s | 1321.32 |
| $\mathrm{K}^{0}$ | ds | 497.672 | $\mathrm{n}^{0}$ | udd | 939.56563 | $\Xi^{+}$ | $\overline{\mathrm{d}} \overline{\mathrm{s}} \bar{s}$ | 1321.32 |
| $\mathrm{K}^{+}$ | us | 493.677 | $\overline{\mathrm{n}}^{0}$ | $\overline{\mathrm{u}} \overline{\mathrm{d}} \overline{\mathrm{d}}$ | 939.56563 | $\Omega^{-}$ | S S S | 1672.45 |
| $\mathrm{K}^{-}$ | sū | 493.677 | $\Lambda$ | uds | 1115.684 | $\Omega^{+}$ | S S S | 1672.45 |
| $\mathrm{D}^{+}$ | c $\bar{d}$ | 1869.4 | $\bar{\Lambda}$ | $\overline{\mathrm{u}} \overline{\mathrm{d}} \overline{\mathrm{s}}$ | 1115.684 | $\Lambda_{c}^{+}$ | udc | 2285.1 |
| $\mathrm{D}^{-}$ | $\mathrm{d} \overline{\mathrm{c}}$ | 1869.4 | $\Sigma \Sigma^{+}$ | u us | 1189.37 | $\Delta^{2-}$ | $\bar{u} \bar{u} \bar{u}$ | 1232.0 |
| $\mathrm{D}^{0}$ | $\mathrm{c} \overline{\mathrm{u}}$ | 1864.6 | $\overline{\Sigma^{-}}$ | $\overline{\mathrm{u}} \mathrm{u} \overline{\mathrm{s}}$ | 1189.37 | $\Delta^{2+}$ | u u u | 1232.0 |
| $\mathrm{D}^{0}$ | uc | 1864.6 | $\Sigma^{0}$ | uds | 1192.55 | $\Delta^{+}$ | u ud | 1232.0 |
| $\mathrm{F}^{+}$ | cs | 1969.0 | $\Sigma^{0}$ | $\overline{\mathrm{u}} \overline{\mathrm{d}} \overline{\mathrm{s}}$ | 1192.55 | $\Delta^{0}$ | udd | 1232.0 |
| $\mathrm{F}^{-}$ | sc | 1969.0 | $\Sigma^{-}$ | d d s | 1197.436 | $\Delta^{-}$ | d d d | 1232.0 |

Each quark can exist in two spin states. So mesons are bosons with spin 0 or 1 in their ground state, while baryons are fermions with spin $\frac{1}{2}$ or $\frac{3}{2}$. There exist excited states with higher internal $L$. Neutrino's have a helicity of $-\frac{1}{2}$ while antineutrino's have only $+\frac{1}{2}$ as possible value.

The quantum numbers are subject to conservation laws. These can be derived from symmetries in the Lagrange density: continuous symmetries give rise to additive conservation laws, discrete symmetries result in multiplicative conservation laws.
Geometrical conservation laws are invariant under Lorentz transformations and the CPT-operation. These are:

1. Mass/energy because the laws of nature are invariant for translations in time.
2. Momentum because the laws of nature are invariant for translations in space.
3. Angular momentum because the laws of nature are invariant for rotations.

Dynamical conservation laws are invariant under the CPT-operation. These are:

1. Electrical charge because the Maxwell equations are invariant under gauge transformations.
2. Colour charge is conserved.
3. Isospin because QCD is invariant for rotations in T-space.
4. Baryon number and lepton number are conserved but not under a possible $\mathrm{SU}(5)$ symmetry of the laws of nature.
5. Quarks type is only conserved under the colour interaction.
6. Parity is conserved except for weak interactions.

The elementary particles can be classified into three families:

|  | leptons | quarks | antileptons | antiquarks |
| :---: | :---: | :---: | :---: | :---: |
| 1st generation | $\mathrm{e}^{-}$ | d | $\mathrm{e}^{+}$ | $\overline{\mathrm{d}}$ |
|  | $\nu_{\mathrm{e}}$ | u | $\bar{\nu}_{\mathrm{e}}$ | $\overline{\mathrm{u}}$ |
| 2nd generation | $\mu^{-}$ | s | $\mu^{+}$ | $\overline{\mathrm{s}}$ |
|  | $\nu_{\mu}$ | c | $\bar{\nu}_{\mu}$ | $\overline{\mathrm{c}}$ |
| 3rd generation | $\tau^{-}$ | b | $\tau^{+}$ | $\overline{\mathrm{b}}$ |
|  | $\nu_{\tau}$ | t | $\bar{\nu}_{\tau}$ | $\overline{\mathrm{t}}$ |

Quarks exist in three colours but because they are confined these colours cannot be seen directly. The color force does not decrease with distance. The potential energy will become high enough to create a quarkantiquark pair when it is tried to disjoin an (anti)quark from a hadron. This will result in two hadrons and not in free quarks.

### 15.12 $P$ and CP-violation

It is found that the weak interaction violates P-symmetry, and even CP-symmetry is not conserved. Some processes which violate P symmetry but conserve the combination CP are:

1. $\mu$-decay: $\mu^{-} \rightarrow \mathrm{e}^{-}+\nu_{\mu}+\bar{\nu}_{\mathrm{e}}$. Left-handed electrons appear more than $1000 \times$ as much as right-handed ones.
2. $\beta$-decay of spin-polarized ${ }^{60} \mathrm{Co}:{ }^{60} \mathrm{Co} \rightarrow{ }^{60} \mathrm{Ni}+\mathrm{e}^{-}+\bar{\nu}_{\mathrm{e}}$. More electrons with a spin parallel to the Co than with a spin antiparallel are created: (parallel-antiparallel)/(total) $=20 \%$.
3. There is no connection with the neutrino: the decay of the $\Lambda$ particle through: $\Lambda \rightarrow \mathrm{p}^{+}+\pi^{-}$and $\Lambda \rightarrow \mathrm{n}^{0}+\pi^{0}$ has also these properties.
The CP-symmetry was found to be violated by the decay of neutral Kaons. These are the lowest possible states with a s-quark so they can decay only weakly. The following holds: $\mathrm{C}\left|\mathrm{K}^{0}\right\rangle=\eta\left|\overline{\mathrm{K}^{0}}\right\rangle$ where $\eta$ is a phase factor. Further holds $\mathrm{P}\left|\mathrm{K}^{0}\right\rangle=-\left|\mathrm{K}^{0}\right\rangle$ because $\mathrm{K}^{0}$ and $\overline{\mathrm{K}^{0}}$ have an intrinsic parity of -1 . From this follows that $\mathrm{K}^{0}$ and $\overline{\mathrm{K}^{0}}$ are not eigenvalues of $\mathrm{CP}: \mathrm{CP}\left|\mathrm{K}^{0}\right\rangle=\left|\overline{\mathrm{K}^{0}}\right\rangle$. The linear combinations

$$
\left|\mathrm{K}_{1}^{0}\right\rangle:=\frac{1}{2} \sqrt{2}\left(\left|\mathrm{~K}^{0}\right\rangle+\left|\overline{\mathrm{K}^{0}}\right\rangle\right) \text { and }\left|\mathrm{K}_{2}^{0}\right\rangle:=\frac{1}{2} \sqrt{2}\left(\left|\mathrm{~K}^{0}\right\rangle-\left|\overline{\mathrm{K}^{0}}\right\rangle\right)
$$

are eigenstates of CP : $\mathrm{CP}\left|\mathrm{K}_{1}^{0}\right\rangle=+\left|\mathrm{K}_{1}^{0}\right\rangle$ and $\mathrm{CP}\left|\mathrm{K}_{2}^{0}\right\rangle=-\left|\mathrm{K}_{2}^{0}\right\rangle$. A base of $\mathrm{K}_{1}^{0}$ and $\mathrm{K}_{2}^{0}$ is practical while describing weak interactions. For colour interactions a base of $\mathrm{K}^{0}$ and $\overline{\mathrm{K}^{0}}$ is practical because then the number u -number $\overline{\mathrm{u}}$ is constant. The expansion postulate must be used for weak decays:

$$
\left|\mathrm{K}^{0}\right\rangle=\frac{1}{2}\left(\left\langle\mathrm{~K}_{1}^{0} \mid \mathrm{K}^{0}\right\rangle+\left\langle\mathrm{K}_{2}^{0} \mid \mathrm{K}^{0}\right\rangle\right)
$$

The probability to find a final state with $\mathrm{CP}=-1$ is $\frac{1}{2}\left|\left\langle\mathrm{~K}_{2}^{0} \mid \mathrm{K}^{0}\right\rangle\right|^{2}$, the probability of $\mathrm{CP}=+1$ decay is $\frac{1}{2}\left|\left\langle\mathrm{~K}_{1}^{0} \mid \mathrm{K}^{0}\right\rangle\right|^{2}$.
The relation between the mass eigenvalues of the quarks (unaccented) and the fields arising in the weak currents (accented) is $\left(u^{\prime}, c^{\prime}, t^{\prime}\right)=(u, c, t)$, and:

$$
\left.\begin{array}{rl}
\left(\begin{array}{c}
d^{\prime} \\
s^{\prime} \\
b^{\prime}
\end{array}\right)= & \left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & \cos \theta_{2} & \sin \theta_{2} \\
0 & -\sin \theta_{2} & \cos \theta_{2}
\end{array}\right)\left(\begin{array}{lll}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & \mathrm{e}^{i \delta}
\end{array}\right)\left(\begin{array}{cc}
\cos \theta_{1} & \sin \theta_{1} \\
0 \\
-\sin \theta_{1} & \cos \theta_{1} \\
0 \\
0 & 0
\end{array}\right) 1
\end{array}\right)
$$

$\theta_{1} \equiv \theta_{\mathrm{C}}$ is the Cabibbo angle: $\sin \left(\theta_{\mathrm{C}}\right) \approx 0.23 \pm 0.01$.

### 15.13 The standard model

When one wants to make the Lagrange density which describes a field invariant for local gauge transformations from a certain group, one has to perform the transformation

$$
\frac{\partial}{\partial x_{\mu}} \rightarrow \frac{D}{D x_{\mu}}=\frac{\partial}{\partial x_{\mu}}-i \frac{g}{\hbar} L_{k} A_{\mu}^{k}
$$

Here the $L_{k}$ are the generators of the gauge group (the "charges") and the $A_{\mu}^{k}$ are the gauge fields. $g$ is the matching coupling constant. The Lagrange density for a scalar field becomes:

$$
\mathcal{L}=-\frac{1}{2}\left(D_{\mu} \Phi^{*} D^{\mu} \Phi+M^{2} \Phi^{*} \Phi\right)-\frac{1}{4} F_{\mu \nu}^{a} F_{a}^{\mu \nu}
$$

and the field tensors are given by: $F_{\mu \nu}^{a}=\partial_{\mu} A_{\nu}^{a}-\partial_{\nu} A_{\mu}^{a}+g c_{l m}^{a} A_{\mu}^{l} A_{\nu}^{m}$.

### 15.13.1 The electroweak theory

The electroweak interaction arises from the necessity to keep the Lagrange density invariant for local gauge transformations of the group $\mathrm{SU}(2) \otimes \mathrm{U}(1)$. Right- and left-handed spin states are treated different because the weak interaction does not conserve parity. If a fifth Dirac matrix is defined by:

$$
\gamma_{5}:=\gamma_{1} \gamma_{2} \gamma_{3} \gamma_{4}=-\left(\begin{array}{cccc}
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1 \\
1 & 0 & 0 & 0 \\
0 & 1 & 0 & 0
\end{array}\right)
$$

the left- and right- handed solutions of the Dirac equation for neutrino's are given by:

$$
\psi_{\mathrm{L}}=\frac{1}{2}\left(1+\gamma_{5}\right) \psi \quad \text { and } \quad \psi_{\mathrm{R}}=\frac{1}{2}\left(1-\gamma_{5}\right) \psi
$$

It appears that neutrino's are always left-handed while antineutrino's are always right-handed. The hypercharge $Y$, for quarks given by $Y=\mathrm{B}+\mathrm{S}+\mathrm{C}+\mathrm{B}^{*}+\mathrm{T}^{\prime}$, is defined by:

$$
Q=\frac{1}{2} Y+T_{3}
$$

so $\left[Y, T_{k}\right]=0$. The group $\mathrm{U}(1)_{Y} \otimes \mathrm{SU}(2)_{T}$ is taken as symmetry group for the electroweak interaction because the generators of this group commute. The multiplets are classified as follows:

|  | $\mathrm{e}_{\mathrm{R}}^{-}$ | $\nu_{\mathrm{eL}} \mathrm{e}_{\mathrm{L}}^{-}$ | $\mathrm{u}_{\mathrm{L}} \mathrm{d}_{\mathrm{L}}^{\prime}$ | $\mathrm{u}_{\mathrm{R}}$ | $\mathrm{d}_{\mathrm{R}}$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $T$ | 0 | $\frac{1}{2}$ | $\frac{1}{2}$ | 0 | 0 |
| $T_{3}$ | 0 | $\frac{1}{2}-\frac{1}{2}$ | $\frac{1}{2}-\frac{1}{2}$ | 0 | 0 |
| $Y$ | -2 | -1 | $\frac{1}{3}$ | $\frac{4}{3}$ | $-\frac{2}{3}$ |

Now, 1 field $B_{\mu}(x)$ is connected with gauge group $\mathrm{U}(1)$ and 3 gauge fields $\vec{A}_{\mu}(x)$ are connected with $\mathrm{SU}(2)$. The total Lagrange density (minus the fieldterms) for the electron-fermion field now becomes:

$$
\begin{aligned}
\mathcal{L}_{0, \mathrm{EW}}= & -\left(\overline{\psi_{\nu \mathrm{e}, \mathrm{~L}}}, \overline{\psi_{\mathrm{eL}}}\right) \gamma^{\mu}\left(\partial_{\mu}-i \frac{g}{\hbar} \vec{A}_{\mu} \cdot\left(\frac{1}{2} \vec{\sigma}\right)-\frac{1}{2} i \frac{g^{\prime}}{\hbar} B_{\mu} \cdot(-1)\right)\binom{\psi_{\nu \mathrm{e}, \mathrm{~L}}}{\psi_{\mathrm{eL}}}- \\
& \overline{\psi_{\mathrm{eR}}} \gamma^{\mu}\left(\partial_{\mu}-\frac{1}{2} i \frac{g^{\prime}}{\hbar}(-2) B_{\mu}\right) \psi_{\mathrm{eR}}
\end{aligned}
$$

Here, $\frac{1}{2} \vec{\sigma}$ are the generators of $T$ and -1 and -2 the generators of $Y$.

### 15.13.2 Spontaneous symmetry breaking: the Higgs mechanism

All leptons are massless in the equations above. Their mass is probably generated by spontaneous symmetry breaking. This means that the dynamic equations which describe the system have a symmetry which the ground state does not have. It is assumed that there exists an isospin-doublet of scalar fields $\Phi$ with electrical charges +1 and 0 and potential $V(\Phi)=-\mu^{2} \Phi^{*} \Phi+\lambda\left(\Phi^{*} \Phi\right)^{2}$. Their antiparticles have charges -1 and 0 . The extra terms in $\mathcal{L}$ arising from these fields, $\mathcal{L}_{H}=\left(D_{L \mu} \Phi\right)^{*}\left(D_{L}^{\mu} \Phi\right)-V(\Phi)$, are globally $\mathrm{U}(1) \otimes \mathrm{SU}(2)$ symmetric. Hence the state with the lowest energy corresponds with the state $\Phi^{*}(x) \Phi(x)=v=\mu^{2} / 2 \lambda=$ constant. The field can be written (were $\omega^{ \pm}$and $z$ are Nambu-Goldstone bosons which can be transformed away, and $m_{\phi}=\mu \sqrt{2}$ ) as:

$$
\Phi=\binom{\Phi^{+}}{\Phi^{0}}=\binom{i \omega^{+}}{(v+\phi-i z) / \sqrt{2}} \text { and }\langle 0| \Phi|0\rangle=\binom{0}{v / \sqrt{2}}
$$

Because this expectation value $\neq 0$ the $S U(2)$ symmetry is broken but the $U(1)$ symmetry is not. When the gauge fields in the resulting Lagrange density are separated one obtains:

$$
\begin{aligned}
W_{\mu}^{-} & =\frac{1}{2} \sqrt{2}\left(A_{\mu}^{1}+i A_{\mu}^{2}\right), W_{\mu}^{+}=\frac{1}{2} \sqrt{2}\left(A_{\mu}^{1}-i A_{\mu}^{2}\right) \\
Z_{\mu} & =\frac{g A_{\mu}^{3}-g^{\prime} B_{\mu}}{\sqrt{g^{2}+g^{\prime 2}}} \equiv A_{\mu}^{3} \cos \left(\theta_{\mathrm{W}}\right)-B_{\mu} \sin \left(\theta_{\mathrm{W}}\right) \\
A_{\mu} & =\frac{g^{\prime} A_{\mu}^{3}+g B_{\mu}}{\sqrt{g^{2}+g^{\prime 2}}} \equiv A_{\mu}^{3} \sin \left(\theta_{\mathrm{W}}\right)+B_{\mu} \cos \left(\theta_{\mathrm{W}}\right)
\end{aligned}
$$

where $\theta_{\mathrm{W}}$ is called the Weinberg angle. For this angle holds: $\sin ^{2}\left(\theta_{\mathrm{W}}\right)=0.255 \pm 0.010$. Relations for the masses of the field quanta can be obtained from the remaining terms: $M_{W}=\frac{1}{2} v g$ and $M_{Z}=\frac{1}{2} v \sqrt{g^{2}+g^{\prime 2}}$, and for the elementary charge holds: $e=\frac{g g^{\prime}}{\sqrt{g^{2}+g^{\prime 2}}}=g^{\prime} \cos \left(\theta_{\mathrm{W}}\right)=g \sin \left(\theta_{\mathrm{W}}\right)$
Experimentally it is found that $M_{W}=80.022 \pm 0.26 \mathrm{GeV} / \mathrm{c}^{2}$ and $M_{Z}=91.187 \pm 0.007 \mathrm{GeV} / \mathrm{c}^{2}$. According to the weak theory this should be: $M_{W}=83.0 \pm 0.24 \mathrm{GeV} / \mathrm{c}^{2}$ and $M_{Z}=93.8 \pm 2.0 \mathrm{GeV} / \mathrm{c}^{2}$.

### 15.13.3 Quantumchromodynamics

Coloured particles interact because the Lagrange density is invariant for the transformations of the group $\mathrm{SU}(3)$ of the colour interaction. A distinction can be made between two types of particles:

1. "White" particles: they have no colour charge, the generator $\vec{T}=0$.
2. "Coloured" particles: the generators $\vec{T}$ are $83 \times 3$ matrices. There exist three colours and three anticolours.
The Lagrange density for coloured particles is given by

$$
\mathcal{L}_{\mathrm{QCD}}=i \sum_{k} \overline{\Psi_{k}} \gamma^{\mu} D_{\mu} \Psi_{k}+\sum_{k, l} \overline{\Psi_{k}} M_{k l} \Psi_{l}-\frac{1}{4} F_{\mu \nu}^{a} F_{a}^{\mu \nu}
$$

The gluons remain massless because this Lagrange density does not contain spinless particles. Because leftand right- handed quarks now belong to the same multiplet a mass term can be introduced. This term can be brought in the form $M_{k l}=m_{k} \delta_{k l}$.

### 15.14 Path integrals

The development in time of a quantum mechanical system can, besides with Schrödingers equation, also be described by a path integral (Feynman):

$$
\psi\left(x^{\prime}, t^{\prime}\right)=\int F\left(x^{\prime}, t^{\prime}, x, t\right) \psi(x, t) d x
$$

in which $F\left(x^{\prime}, t^{\prime}, x, t\right)$ is the amplitude of probability to find a system on time $t^{\prime}$ in $x^{\prime}$ if it was in $x$ on time $t$. Then,

$$
F\left(x^{\prime}, t^{\prime}, x, t\right)=\int \exp \left(\frac{i S[x]}{\hbar}\right) d[x]
$$

where $S[x]$ is an action-integral: $S[x]=\int L(x, \dot{x}, t) d t$. The notation $d[x]$ means that the integral has to be taken over all possible paths $[x]$ :

$$
\int d[x]:=\lim _{n \rightarrow \infty} \frac{1}{N} \prod_{n}\left\{\int_{-\infty}^{\infty} d x\left(t_{n}\right)\right\}
$$

in which $N$ is a normalization constant. To each path is assigned a probability amplitude $\exp (i S / \hbar)$. The classical limit can be found by taking $\delta S=0$ : the average of the exponent vanishes, except where it is stationary. In quantum fieldtheory, the probability of the transition of a fieldoperator $\Phi(\vec{x},-\infty)$ to $\Phi^{\prime}(\vec{x}, \infty)$ is given by

$$
F\left(\Phi^{\prime}(\vec{x}, \infty), \Phi(\vec{x},-\infty)\right)=\int \exp \left(\frac{i S[\Phi]}{\hbar}\right) d[\Phi]
$$

with the action-integral

$$
S[\Phi]=\int_{\Omega} \mathcal{L}\left(\Phi, \partial_{\nu} \Phi\right) d^{4} x
$$

### 15.15 Unification and quantum gravity

The strength of the forces varies with energy and the reciprocal coupling constants approach each other with increasing energy. The $S U(5)$ model predicts complete unification of the electromagnetical, weak and colour forces at $10^{15} \mathrm{GeV}$. It also predicts 12 extra X bosons which couple leptons and quarks and are i.g. responsible for proton decay, with dominant channel $\mathrm{p}^{+} \rightarrow \pi^{0}+\mathrm{e}^{+}$, with an average lifetime of the proton of $10^{31}$ year. This model has been experimentally falsified.

Supersymmetric models assume a symmetry between bosons and fermions and predict partners for the currently known particles with a spin which differs $\frac{1}{2}$. The supersymmetric $\operatorname{SU}(5)$ model predicts unification at $10^{16} \mathrm{GeV}$ and an average lifetime of the proton of $10^{33}$ year. The dominant decay channels in this theory are $\mathrm{p}^{+} \rightarrow \mathrm{K}^{+}+\bar{\nu}_{\mu}$ and $\mathrm{p}^{+} \rightarrow \mathrm{K}^{0}+\mu^{+}$.

Quantum gravity plays only a role in particle interactions at the Planck dimensions, where $\lambda_{\mathrm{C}} \approx R_{\mathrm{S}}: m_{\mathrm{Pl}}=$ $\sqrt{h c / G}=3 \cdot 10^{19} \mathrm{GeV}, t_{\mathrm{Pl}}=h / m_{\mathrm{Pl}} c^{2}=\sqrt{h G / c^{5}}=10^{-43} \mathrm{sec}$ and $r_{\mathrm{Pl}}=c t_{\mathrm{Pl}} \approx 10^{-35} \mathrm{~m}$.

## Chapter 16

## Astrophysics

### 16.1 Determination of distances

The parallax is mostly used to determine distances in nearby space. The parallax is the angular difference between two measurements of the position of the object from different view-points. If the annual parallax is given by $p$, the distance $R$ of the object is given by $R=a / \sin (p)$, in which $a$ is the radius of the Earth's orbit. The clusterparallax is used to determine the distance of a group of stars by using their motion w.r.t. a fixed background. The tangential velocity $v_{\mathrm{t}}$ and the radial velocity $v_{\mathrm{r}}$ of the stars along the sky are given by

$$
v_{\mathrm{r}}=V \cos (\theta) \quad, \quad v_{\mathrm{t}}=V \sin (\theta)=\omega R
$$

where $\theta$ is the angle between the star and the point of convergence and $\hat{R}$ the distance in pc. This results, with $v_{\mathrm{t}}=v_{\mathrm{r}} \tan (\theta)$, in:

$$
R=\frac{v_{\mathrm{r}} \tan (\theta)}{\omega} \Rightarrow \hat{R}=\frac{1^{\prime \prime}}{p}
$$

where $p$ is the parallax in arc seconds. The parallax is then given by

$$
p=\frac{4.74 \mu}{v_{\mathrm{r}} \tan (\theta)}
$$


with $\mu$ de proper motion of the star in "/yr. A method to determine the distance of objects which are somewhat further away, like galaxies and star clusters, uses the period-Brightness relation for Cepheids. This relation is shown in the above figure for different types of stars.

### 16.2 Brightness and magnitudes

The brightness is the total radiated energy per unit of time. Earth receives $s_{0}=1.374 \mathrm{~kW} / \mathrm{m}^{2}$ from the Sun. Hence, the brightness of the Sun is given by $L_{\odot}=4 \pi r^{2} s_{0}=3.82 \cdot 10^{26} \mathrm{~W}$. It is also given by:

$$
L_{\odot}=4 \pi R_{\odot}^{2} \int_{0}^{\infty} \pi F_{\nu} d \nu
$$

where $\pi F_{\nu}$ is the monochromatic radiation flux. At the position of an observer this is $\pi f_{\nu}$, with $f_{\nu}=(R / r)^{2} F_{\nu}$ if absorption is ignored. If $A_{\nu}$ is the fraction of the flux which reaches Earth's surface, the transmission factor is given by $R_{\nu}$ and the surface of the detector is given by $\pi a^{2}$, then the apparent brightness $b$ is given by:

$$
b=\pi a^{2} \int_{0}^{\infty} f_{\nu} A_{\nu} R_{\nu} d \nu
$$

The magnitude $m$ is defined by:

$$
\frac{b_{1}}{b_{2}}=(100)^{\frac{1}{5}\left(m_{2}-m_{1}\right)}=(2.512)^{m_{2}-m_{1}}
$$

because the human eye perceives lightintensities logaritmical. From this follows that $m_{2}-m_{1}=2.5 \cdot{ }^{10}$ $\log \left(b_{1} / b_{2}\right)$, or: $m=-2.5 \cdot{ }^{10} \log (b)+C$. The apparent brightness of a star if this star would be at a distance of 10 pc is called the absolute brightness $B: B / b=(\hat{r} / 10)^{2}$. The absolute magnitude is then given by $M=-2.5 \cdot{ }^{10} \log (B)+C$, or: $M=5+m-5 \cdot{ }^{10} \log (\hat{r})$. When an interstellar absorption of $10^{-4} / \mathrm{pc}$ is taken into account one finds:

$$
M=\left(m-4 \cdot 10^{-4} \hat{r}\right)+5-5 \cdot{ }^{10} \log (\hat{r})
$$

If a detector detects all radiation emitted by a source one would measure the absolute bolometric magnitude. If the bolometric correction $B C$ is given by

$$
B C=2.5 \cdot{ }^{10} \log \left(\frac{\text { Energy flux received }}{\text { Energy flux detected }}\right)=2.5 \cdot{ }^{10} \log \left(\frac{\int f_{\nu} d \nu}{\int f_{\nu} A_{\nu} R_{\nu} d \nu}\right)
$$

holds: $M_{b}=M_{V}-B C$ where $M_{V}$ is the visual magnitude. Further holds

$$
M_{b}=-2.5 \cdot{ }^{10} \log \left(\frac{L}{L_{\odot}}\right)+4.72
$$

### 16.3 Radiation and stellar atmospheres

The radiation energy passing through a surface $d A$ is $d E=I_{\nu}(\theta, \varphi) \cos (\theta) d \nu d \Omega d A d t$, where $I_{\mu}$ is the monochromatical intensity $\left[\mathrm{Wm}^{-2} \mathrm{sr}^{-1} \mathrm{~Hz}^{-1}\right]$. When there is no absorption the quantity $I_{\nu}$ is independent of the distance to the source. Planck's law holds for a black body:

$$
I_{\nu}(T) \equiv B_{\nu}(T)=\frac{c}{4 \pi} w_{\nu}(T)=\frac{2 h \nu^{3}}{c^{2}} \frac{1}{\exp (h \nu / k T)-1}
$$

The radiation transport through a layer can then be written as:

$$
\frac{d I_{\nu}}{d s}=-I_{\nu} \kappa_{\nu}+j_{\nu}
$$

Here, $j_{\nu}$ is the coefficient of emission and $\kappa_{\nu}$ the coefficient of absorption. $\int d s$ is the thickness of the layer. The optical thickness $\tau_{\nu}$ of the layer is given by $\tau_{\nu}=\int \kappa_{\nu} d s$. The layer is optically thin if $\tau_{\nu} \ll 1$, the layer is optically thick if $\tau_{\nu} \gg 1$. For a stellar atmosphere in LTE holds: $j_{\nu}=\kappa_{\nu} B_{\nu}(T)$. Then also holds:

$$
I_{\nu}(s)=I_{\nu}(0) \mathrm{e}^{-\tau_{\nu}}+B_{\nu}(T)\left(1-\mathrm{e}^{-\tau_{\nu}}\right)
$$

### 16.4 Composition and evolution of stars

The structure of a star is described by the following equations:

$$
\begin{aligned}
\frac{d M(r)}{d r} & =4 \pi \varrho(r) r^{2} \\
\frac{d p(r)}{d r} & =-\frac{G M(r) \varrho(r)}{r^{2}} \\
\frac{L(r)}{d r} & =4 \pi \varrho(r) \varepsilon(r) r^{2} \\
\left(\frac{d T(r)}{d r}\right)_{\mathrm{rad}} & =-\frac{3}{4} \frac{L(r)}{4 \pi r^{2}} \frac{\kappa(r)}{4 \sigma T^{3}(r)}, \text { (Eddington), or } \\
\left(\frac{d T(r)}{d r}\right)_{\mathrm{conv}} & =\frac{T(r)}{p(r)} \frac{\gamma-1}{\gamma} \frac{d p(r)}{d r}, \text { (convective energy transport) }
\end{aligned}
$$

Further, for stars of the solar type, the composing plasma can be described as an ideal gas:

$$
p(r)=\frac{\varrho(r) k T(r)}{\mu m_{\mathrm{H}}}
$$

where $\mu$ is the average molecular mass, usually well approximated by:

$$
\mu=\frac{\varrho}{n m_{\mathrm{H}}}=\frac{1}{2 X+\frac{3}{4} Y+\frac{1}{2} Z}
$$

where $X$ is the mass fraction of $\mathrm{H}, Y$ the mass fraction of He and $Z$ the mass fraction of the other elements. Further holds:

$$
\kappa(r)=f(\varrho(r), T(r), \text { composition }) \text { and } \varepsilon(r)=g(\varrho(r), T(r), \text { composition })
$$

Convection will occur when the star meets the Schwartzschild criterium:

$$
\left(\frac{d T}{d r}\right)_{\mathrm{conv}}<\left(\frac{d T}{d r}\right)_{\mathrm{rad}}
$$

Otherwise the energy transfer takes place by radiation. For stars in quasi-hydrostatic equilibrium hold the approximations $r=\frac{1}{2} R, M(r)=\frac{1}{2} M, d M / d r=M / R, \kappa \sim \varrho$ and $\varepsilon \sim \varrho T^{\mu}$ (this last assumption is only valid for stars on the main sequence). For pp-chains holds $\mu \approx 5$ and for the CNO chains holds $\mu=12$ tot 18 . It can be derived that $L \sim M^{3}$ : the mass-brightness relation. Further holds: $L \sim R^{4} \sim T_{\text {eff }}^{8}$. This results in the equation for the main sequence in the Hertzsprung-Russel diagram:

$$
{ }^{10} \log (L)=8 \cdot{ }^{10} \log \left(T_{\text {eff }}\right)+\text { constant }
$$

### 16.5 Energy production in stars

The net reaction from which most stars gain their energy is: $4^{1} \mathrm{H} \rightarrow{ }^{4} \mathrm{He}+2 \mathrm{e}^{+}+2 \nu_{\mathrm{e}}+\gamma$.
This reaction produces 26.72 MeV . Two reaction chains are responsible for this reaction. The slowest, speedlimiting reaction is shown in boldface. The energy between brackets is the energy carried away by the neutrino.

1. The proton-proton chain can be divided into two subchains:
${ }^{1} \mathbf{H}+\mathbf{p}^{+} \rightarrow{ }^{2} \mathbf{D}+\mathbf{e}^{+}+\nu_{\mathbf{e}}$, and then ${ }^{2} \mathrm{D}+\mathrm{p} \rightarrow{ }^{3} \mathrm{He}+\gamma$.
I. pp1: ${ }^{3} \mathrm{He}+{ }^{3} \mathrm{He} \rightarrow 2 \mathrm{p}^{+}+{ }^{4} \mathrm{He}$. There is $26.21+(0.51) \mathrm{MeV}$ released.
II. pp2: ${ }^{3} \mathrm{He}+\alpha \rightarrow{ }^{7} \mathrm{Be}+\gamma$
i. ${ }^{7} \mathrm{Be}+\mathrm{e}^{-} \rightarrow{ }^{7} \mathrm{Li}+\nu$, then ${ }^{7} \mathrm{Li}+\mathrm{p}^{+} \rightarrow 2^{4} \mathrm{He}+\gamma .25 .92+(0.80) \mathrm{MeV}$.
ii. ${ }^{7} \mathrm{Be}+\mathrm{p}^{+} \rightarrow{ }^{8} \mathrm{~B}+\gamma$, then ${ }^{8} \mathrm{~B}+\mathrm{e}^{+} \rightarrow 2^{4} \mathrm{He}+\bar{\nu} .19 .5+(7.2) \mathrm{MeV}$.

Both ${ }^{7} \mathrm{Be}$ chains become more important with raising $T$.
2. The CNO cycle. The first chain releases $25.03+(1.69) \mathrm{MeV}$, the second $24.74+(1.98) \mathrm{MeV}$. The reactions are shown below.

$$
\begin{aligned}
& \rightarrow \quad{ }^{15} \mathrm{~N}+\mathrm{p}^{+} \underset{\downarrow}{\longrightarrow} \alpha+{ }^{12} \mathrm{C} \quad{ }^{15} \mathrm{~N}+\mathrm{p}^{+} \rightarrow{ }^{16} \mathrm{O}+\gamma \\
& { }^{12} \mathrm{C}+\mathrm{p}^{+} \rightarrow{ }^{13} \mathrm{~N}+\gamma \quad{ }^{16} \mathrm{O}+\mathrm{p}^{+} \xrightarrow{\downarrow}{ }^{17} \mathrm{~F}+\gamma \\
& { }^{14} \mathbf{N}+\mathrm{p}^{+} \xrightarrow{\uparrow}{ }^{\mathbf{1 5}} \mathrm{O}+\gamma \\
& { }^{13} \mathrm{~N} \rightarrow \stackrel{\downarrow}{13} \mathrm{C}+\mathrm{e}^{+}+\nu \\
& { }^{17} \mathrm{~F} \rightarrow \stackrel{\downarrow}{17} \mathrm{O}+\mathrm{e}^{+}+\nu \\
& \leftarrow \quad{ }^{13} \mathrm{C}+\mathrm{p}^{+} \xrightarrow{\downarrow}{ }^{14} \mathrm{~N}+\gamma \\
& { }^{17} \mathrm{O}+\mathrm{p}^{+} \xrightarrow{\downarrow} \alpha+{ }^{14} \mathrm{~N}
\end{aligned}
$$

## The $\nabla$-operator

In cartesian coordinates $(x, y, z)$ holds:

$$
\begin{gathered}
\vec{\nabla}=\frac{\partial}{\partial x} \vec{e}_{x}+\frac{\partial}{\partial y} \vec{e}_{y}+\frac{\partial}{\partial z} \vec{e}_{z}, \operatorname{grad} f=\vec{\nabla} f=\frac{\partial f}{\partial x} \vec{e}_{x}+\frac{\partial f}{\partial y} \vec{e}_{y}+\frac{\partial f}{\partial z} \vec{e}_{z} \\
\operatorname{div} \vec{a}=\vec{\nabla} \cdot \vec{a}=\frac{\partial a_{x}}{\partial x}+\frac{\partial a_{y}}{\partial y}+\frac{\partial a_{z}}{\partial z}, \nabla^{2} f=\frac{\partial^{2} f}{\partial x^{2}}+\frac{\partial^{2} f}{\partial y^{2}}+\frac{\partial^{2} f}{\partial z^{2}} \\
\operatorname{rot} \vec{a}=\vec{\nabla} \times \vec{a}=\left(\frac{\partial a_{z}}{\partial y}-\frac{\partial a_{y}}{\partial z}\right) \vec{e}_{x}+\left(\frac{\partial a_{x}}{\partial z}-\frac{\partial a_{z}}{\partial x}\right) \vec{e}_{y}+\left(\frac{\partial a_{y}}{\partial x}-\frac{\partial a_{x}}{\partial y}\right) \vec{e}_{z}
\end{gathered}
$$

In cylinder coordinates $(r, \varphi, z)$ holds:

$$
\begin{gathered}
\vec{\nabla}=\frac{\partial}{\partial r} \vec{e}_{r}+\frac{1}{r} \frac{\partial}{\partial \varphi} \vec{e}_{\varphi}+\frac{\partial}{\partial z} \vec{e}_{z}, \operatorname{grad} f=\frac{\partial f}{\partial r} \vec{e}_{r}+\frac{1}{r} \frac{\partial f}{\partial \varphi} \vec{e}_{\varphi}+\frac{\partial f}{\partial z} \vec{e}_{z} \\
\operatorname{div} \vec{a}=\frac{\partial a_{r}}{\partial r}+\frac{a_{r}}{r}+\frac{1}{r} \frac{\partial a_{\varphi}}{\partial \varphi}+\frac{\partial a_{z}}{\partial z}, \nabla^{2} f=\frac{\partial^{2} f}{\partial r^{2}}+\frac{1}{r} \frac{\partial f}{\partial r}+\frac{1}{r^{2}} \frac{\partial^{2} f}{\partial \varphi^{2}}+\frac{\partial^{2} f}{\partial z^{2}} \\
\operatorname{rot} \vec{a}=\left(\frac{1}{r} \frac{\partial a_{z}}{\partial \varphi}-\frac{\partial a_{\varphi}}{\partial z}\right) \vec{e}_{r}+\left(\frac{\partial a_{r}}{\partial z}-\frac{\partial a_{z}}{\partial r}\right) \vec{e}_{\varphi}+\left(\frac{\partial a_{\varphi}}{\partial r}+\frac{a_{\varphi}}{r}-\frac{1}{r} \frac{\partial a_{r}}{\partial \varphi}\right) \vec{e}_{z}
\end{gathered}
$$

In spherical coordinates $(r, \theta, \varphi)$ holds:

$$
\begin{aligned}
\vec{\nabla}= & \frac{\partial}{\partial r} \vec{e}_{r}+\frac{1}{r} \frac{\partial}{\partial \theta} \vec{e}_{\theta}+\frac{1}{r \sin \theta} \frac{\partial}{\partial \varphi} \vec{e}_{\varphi} \\
\operatorname{grad} f= & \frac{\partial f}{\partial r} \vec{e}_{r}+\frac{1}{r} \frac{\partial f}{\partial \theta} \vec{e}_{\theta}+\frac{1}{r \sin \theta} \frac{\partial f}{\partial \varphi} \vec{e}_{\varphi} \\
\operatorname{div} \vec{a}= & \frac{\partial a_{r}}{\partial r}+\frac{2 a_{r}}{r}+\frac{1}{r} \frac{\partial a_{\theta}}{\partial \theta}+\frac{a_{\theta}}{r \tan \theta}+\frac{1}{r \sin \theta} \frac{\partial a_{\varphi}}{\partial \varphi} \\
\operatorname{rot} \vec{a}= & \left(\frac{1}{r} \frac{\partial a_{\varphi}}{\partial \theta}+\frac{a_{\theta}}{r \tan \theta}-\frac{1}{r \sin \theta} \frac{\partial a_{\theta}}{\partial \varphi}\right) \vec{e}_{r}+\left(\frac{1}{r \sin \theta} \frac{\partial a_{r}}{\partial \varphi}-\frac{\partial a_{\varphi}}{\partial r}-\frac{a_{\varphi}}{r}\right) \vec{e}_{\theta}+ \\
& \left(\frac{\partial a_{\theta}}{\partial r}+\frac{a_{\theta}}{r}-\frac{1}{r} \frac{\partial a_{r}}{\partial \theta}\right) \vec{e}_{\varphi} \\
\nabla^{2} f= & \frac{\partial^{2} f}{\partial r^{2}}+\frac{2}{r} \frac{\partial f}{\partial r}+\frac{1}{r^{2}} \frac{\partial^{2} f}{\partial \theta^{2}}+\frac{1}{r^{2} \tan \theta} \frac{\partial f}{\partial \theta}+\frac{1}{r^{2} \sin ^{2} \theta} \frac{\partial^{2} f}{\partial \varphi^{2}}
\end{aligned}
$$

General orthonormal curvelinear coordinates $(u, v, w)$ can be obtained from cartesian coordinates by the transformation $\vec{x}=\vec{x}(u, v, w)$. The unit vectors are then given by:

$$
\vec{e}_{u}=\frac{1}{h_{1}} \frac{\partial \vec{x}}{\partial u}, \quad \vec{e}_{v}=\frac{1}{h_{2}} \frac{\partial \vec{x}}{\partial v}, \quad \vec{e}_{w}=\frac{1}{h_{3}} \frac{\partial \vec{x}}{\partial w}
$$

where the factors $h_{i}$ set the norm to 1 . Then holds:

$$
\begin{aligned}
\operatorname{grad} f= & \frac{1}{h_{1}} \frac{\partial f}{\partial u} \vec{e}_{u}+\frac{1}{h_{2}} \frac{\partial f}{\partial v} \vec{e}_{v}+\frac{1}{h_{3}} \frac{\partial f}{\partial w} \vec{e}_{w} \\
\operatorname{div} \vec{a}= & \frac{1}{h_{1} h_{2} h_{3}}\left(\frac{\partial}{\partial u}\left(h_{2} h_{3} a_{u}\right)+\frac{\partial}{\partial v}\left(h_{3} h_{1} a_{v}\right)+\frac{\partial}{\partial w}\left(h_{1} h_{2} a_{w}\right)\right) \\
\operatorname{rot} \vec{a}= & \frac{1}{h_{2} h_{3}}\left(\frac{\partial\left(h_{3} a_{w}\right)}{\partial v}-\frac{\partial\left(h_{2} a_{v}\right)}{\partial w}\right) \vec{e}_{u}+\frac{1}{h_{3} h_{1}}\left(\frac{\partial\left(h_{1} a_{u}\right)}{\partial w}-\frac{\partial\left(h_{3} a_{w}\right)}{\partial u}\right) \vec{e}_{v}+ \\
& \frac{1}{h_{1} h_{2}}\left(\frac{\partial\left(h_{2} a_{v}\right)}{\partial u}-\frac{\partial\left(h_{1} a_{u}\right)}{\partial v}\right) \vec{e}_{w} \\
\nabla^{2} f= & \frac{1}{h_{1} h_{2} h_{3}}\left[\frac{\partial}{\partial u}\left(\frac{h_{2} h_{3}}{h_{1}} \frac{\partial f}{\partial u}\right)+\frac{\partial}{\partial v}\left(\frac{h_{3} h_{1}}{h_{2}} \frac{\partial f}{\partial v}\right)+\frac{\partial}{\partial w}\left(\frac{h_{1} h_{2}}{h_{3}} \frac{\partial f}{\partial w}\right)\right]
\end{aligned}
$$

## The SI units

## Basic units

| Quantity | Unit | Sym. |
| :--- | :--- | :--- |
| Length | metre | m |
| Mass | kilogram | kg |
| Time | second | s |
| Therm. temp. | kelvin | K |
| Electr. current | ampere | A |
| Luminous intens. | candela | cd |
| Amount of subst. | mol | mol |

## Extra units

| Plane angle <br> solid angle | radian <br> sterradian | rad <br> sr |
| :--- | :--- | :--- |

Derived units with special names

| Quantity | Unit | Sym. | Derivation |
| :--- | :--- | :--- | :--- |
| Frequency | hertz | Hz | $\mathrm{s}^{-1}$ |
| Force | newton | N | $\mathrm{kg} \cdot \mathrm{m} \cdot \mathrm{s}^{-2}$ |
| Pressure | pascal | Pa | $\mathrm{N} \cdot \mathrm{m}^{-2}$ |
| Energy | joule | J | $\mathrm{N} \cdot \mathrm{m}^{-1}$ |
| Power | watt | W | $\mathrm{J} \cdot \mathrm{s}^{-1}$ |
| Charge | coulomb | C | $\mathrm{A} \cdot \mathrm{s}$ |
| El. Potential | volt | V | $\mathrm{W} \cdot \mathrm{A}^{-1}$ |
| El. Capacitance | farad | F | $\mathrm{C} \cdot \mathrm{V}^{-1}$ |
| El. Resistance | ohm | $\Omega$ | $\mathrm{V} \cdot \mathrm{A}^{-1}$ |
| El. Conductance | siemens | S | $\mathrm{A} \cdot \mathrm{V}^{-1}$ |
| Mag. flux | weber | Wb | $\mathrm{V} \cdot \mathrm{s}$ |
| Mag. flux density | tesla | T | $\mathrm{Wb} \cdot \mathrm{m}^{-2}$ |
| Inductance | henry | H | $\mathrm{Wb} \cdot \mathrm{A}^{-1}$ |
| Luminous flux | lumen | lm | $\mathrm{cd} \cdot \mathrm{sr}^{\text {lin }}$ |
| Illuminance | lux | lx | $\mathrm{lm} \cdot \mathrm{m}^{-2}$ |
| Activity | bequerel | Bq | $\mathrm{s}{ }^{-1}$ |
| Absorbed dose | gray | Gy | $\mathrm{J} \cdot \mathrm{kg}^{-1}$ |
| Dose equivalent | sievert | Sv | $\mathrm{J} \cdot \mathrm{kg}^{-1}$ |

## Prefixes

| yotta | Y | $10^{24}$ | giga | G | $10^{9}$ | deci | d | $10^{-1}$ | pico | p | $10^{-12}$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| zetta | Z | $10^{21}$ | mega | M | $10^{6}$ | centi | c | $10^{-2}$ | femto | f | $10^{-15}$ |
| exa | E | $10^{18}$ | kilo | k | $10^{3}$ | milli | m | $10^{-3}$ | atto | a | $10^{-18}$ |
| peta | P | $10^{15}$ | hecto | h | $10^{2}$ | micro | $\mu$ | $10^{-6}$ | zepto | z | $10^{-21}$ |
| tera | T | $10^{12}$ | deca | da | 10 | nano | n | $10^{-9}$ | yocto | y | $10^{-24}$ |

