## Lecture notes Particle Physics II

## Quantum Chromo Dynamics

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## Preliminary

This section is not part of the lectures, but a small collection of material that should be familiar from special relativity, electrodynamics, quantum mechanics and the lecture series Particle Physics I.

Also included is a summary of group theory, but still very incomplete.

## Units and conversion factors

In particle physics, energy is measured in units of $\mathrm{GeV}=10^{6} \mathrm{eV}$, where $1 \mathrm{eV}=1.6 \times 10^{-19} \mathrm{~J}$ is the change in kinetic energy of an electron when it traverses a potential difference of one volt. From the relation $E^{2}=p^{2} c^{2}+m^{2} c^{4}$ it follows that the units of momentum and mass are $\mathrm{GeV} / c$ and $\mathrm{GeV} / c^{2}$, respectively. The dimension of $\hbar$ is energy $\times$ time so that the unit of time is $\hbar / \mathrm{GeV}$; $\hbar c$ has dimension energy $\times$ length so that length has unit $\hbar c / \mathrm{GeV}$.
One often works in a system of units where $\hbar$ and $c$ have a numerical value of one, so that these constants can be omitted in expressions, as in $E^{2}=p^{2}+m^{2}$. A disadvantage is that the dimensions carried by $\hbar$ (energy $\times$ time) and $c$ (length/time) also disappear but these can always be restored, if necessary, by a dimensional analysis afterward.
Here are some useful conversions.

|  | Conversion | $\hbar=c=1$ units | Natural units |
| :--- | :--- | :---: | :---: |
| Mass | $1 \mathrm{~kg}=5.61 \times 10^{26}$ | GeV | $\mathrm{GeV} / c^{2}$ |
| Length | $1 \mathrm{~m}=5.07 \times 10^{15}$ | $\mathrm{GeV}^{-1}$ | $\hbar c / \mathrm{GeV}$ |
| Time | $1 \mathrm{~s}=1.52 \times 10^{24}$ | $\mathrm{GeV}^{-1}$ | $\hbar / \mathrm{GeV}$ |
| Charge | $e=\sqrt{4 \pi \alpha}$ | dimensionless | $\sqrt{\hbar c}$ |

$$
\begin{aligned}
& 1 \mathrm{TeV}=10^{3} \mathrm{GeV}=10^{6} \mathrm{MeV}=10^{9} \mathrm{KeV}=10^{12} \mathrm{eV} \\
& 1 \mathrm{fm}=10^{-15} \mathrm{~m}=10^{-13} \mathrm{~cm}=5.07 \mathrm{GeV}^{-1} \\
& 1 \mathrm{barn}=10^{-28} \mathrm{~m}^{2}=10^{-24} \mathrm{~cm}^{2} \\
& 1 \mathrm{fm}^{2}=10 \mathrm{mb}=10^{4} \mu \mathrm{~b}=10^{7} \mathrm{nb}=10^{10} \mathrm{pb} \\
& 1 \mathrm{GeV}^{-2}=0.389 \mathrm{mb} \\
& \hbar c=197 \mathrm{MeV} \mathrm{fm} \\
& (\hbar c)^{2}=0.389 \mathrm{GeV}^{2} \mathrm{mb} \\
& \alpha=e^{2} /(4 \pi \hbar c) \approx 1 / 137
\end{aligned}
$$

## Covariant notation ( $c=1$ )

- Contravariant space-time coordinate: $x^{\mu}=\left(x^{0}, x^{1}, x^{2}, x^{3}\right)=(t, \boldsymbol{x})$
- Covariant space-time coordinate: $x_{\mu}=\left(x_{0}, x_{1}, x_{2}, x_{3}\right)=(t,-\boldsymbol{x})$
- Contravariant derivative: $\partial_{\mu} \equiv \partial / \partial x^{\mu}=\left(\partial_{t},+\nabla\right)$
- Covariant derivative: $\quad \partial^{\mu} \equiv \partial / \partial x_{\mu}=\left(\partial_{t},-\nabla\right)$
- Metric tensor: $\quad g_{\mu \nu}=g^{\mu \nu}=\operatorname{diag}(1,-1,-1,-1)$
- Index raising/lowering: $a_{\mu}=g_{\mu \nu} a^{\nu}, \quad a^{\mu}=g^{\mu \nu} a_{\nu}$
- Lorentz boost along $x$-axis: ${ }^{2} x^{\mu}=\Lambda^{\mu}{ }_{\nu} x^{\nu}$

$$
\Lambda_{\nu}^{\mu}=\left(\begin{array}{cccc}
\gamma & -\gamma \beta & 0 & 0 \\
-\gamma \beta & \gamma & 0 & 0 \\
0 & 0 & 1 & 0 \\
0 & 0 & 0 & 1
\end{array}\right) \quad \gamma=\frac{1}{\sqrt{1-\beta^{2}}}
$$

We have also: $x_{\mu}^{\prime}=\Lambda_{\mu}^{\nu} x_{\nu}$ with $\Lambda_{\mu}{ }^{\nu}(\beta)=\Lambda_{\nu}^{\mu}(-\beta) \equiv\left(\Lambda_{\nu}^{\mu}\right)^{-1}$

- Inproduct (Lorentz scalar): $a \cdot b=a^{\mu} b_{\mu}=a^{0} b_{0}-\boldsymbol{a} \cdot \boldsymbol{b}=a_{\mu} b^{\mu}$
- $a^{2}>0$ time-like 4 -vector $\rightarrow$ possible causal connection $a^{2}=0$ light-like 4 -vector
$a^{2}<0$ space-like 4 -vector $\rightarrow$ no causal connection
- 4-momentum: $p^{\mu}=(E, \boldsymbol{p}), \quad p_{\mu}=(E,-\boldsymbol{p})$
- Invariant mass: $p^{2}=p^{\mu} p_{\mu}=p_{\mu} p^{\mu}=E^{2}-\boldsymbol{p}^{2}=m^{2}$
- Particle velocity: $\gamma=E / m, \quad \beta=|\boldsymbol{p}| / E$

[^0]
## Vector calculus

$$
\begin{aligned}
\nabla \times(\nabla \psi) & =\mathbf{0} \\
\nabla \cdot(\nabla \times \boldsymbol{A}) & =0 \\
\nabla \times(\nabla \times \boldsymbol{A}) & =\nabla(\nabla \cdot \boldsymbol{A})-\nabla^{2} \boldsymbol{A} \\
\int_{V} \nabla \cdot \boldsymbol{A} \mathrm{~d} V & =\int_{S} \boldsymbol{A} \cdot \hat{\boldsymbol{n}} \mathrm{~d} S \quad \text { (Divergence theorem) } \\
\int_{V}\left(\phi \nabla^{2} \psi-\psi \nabla^{2} \phi\right) \mathrm{d} V & =\int_{S}(\phi \nabla \psi-\psi \nabla \phi) \cdot \hat{\boldsymbol{n}} d S \quad \text { (Green's theorem) } \\
\int_{S}(\nabla \times \boldsymbol{A}) \cdot \hat{\boldsymbol{n}} \mathrm{d} S & =\oint_{C} \boldsymbol{A} \cdot \mathrm{~d} \boldsymbol{l} \quad \text { (Stokes' theorem) }
\end{aligned}
$$

- In the above, $S$ is a closed surface bounding $V$, with $\hat{\boldsymbol{n}}$ the outward normal unit vector at the surface element $\mathrm{d} S$.
- In Stokes' theorem, the direction of $\hat{\boldsymbol{n}}$ is related by the right-hand rule to the sense of the contour integral around $C$.


## Maxwell's equations in vacuum

- Maxwell's equations

$$
\begin{array}{ll}
\nabla \cdot \boldsymbol{E}=\rho & \nabla \cdot \boldsymbol{B}=0 \\
\nabla \times \boldsymbol{E}+\partial \boldsymbol{B} / \partial t=0 & \nabla \times \boldsymbol{B}-\partial \boldsymbol{E} / \partial t=\boldsymbol{j}
\end{array}
$$

- Continuity equation

$$
\nabla \cdot \boldsymbol{j}=-\frac{\partial \rho}{\partial t}
$$

- The potentials $V$ and $\boldsymbol{A}$ are defined such that the second and third of Maxwell's equations are automatically satisfied

$$
\begin{array}{lll}
\boldsymbol{B}=\nabla \times \boldsymbol{A} & \rightarrow & \nabla \cdot \boldsymbol{B}=0 \\
\boldsymbol{E}=-\partial \boldsymbol{A} / \partial t-\nabla V & \rightarrow & \nabla \times \boldsymbol{E}=-\partial \boldsymbol{B} / \partial t
\end{array}
$$

- Gauge transformations leave the $\boldsymbol{E}$ and $\boldsymbol{B}$ fields invariant

$$
V^{\prime}=V+\frac{\partial \lambda}{\partial t} \quad \text { and } \quad \boldsymbol{A}^{\prime}=\boldsymbol{A}-\nabla \lambda
$$

- Maxwell's equations in 4-vector notation

$$
\begin{array}{ll}
\text { 4-vector potential } & A^{\mu}=(V, \boldsymbol{A}) \\
\text { 4-vector current } & j^{\mu}=(\rho, \boldsymbol{j}) \\
\text { Electromagnetic tensor } & F^{\mu \nu}=\partial^{\mu} A^{\nu}-\partial^{\nu} A^{\mu} \\
\text { Maxwell's equations } & \partial_{\mu} F^{\mu \nu}=j^{\nu} \\
\text { Continuity equation } & \partial_{\mu} j^{\mu}=0 \\
\text { Gauge transformation } & A^{\mu} \rightarrow A^{\mu}+\partial^{\mu} \lambda
\end{array}
$$

- Lorentz gauge and Coulomb condition

Lorentz gauge

$$
\partial_{\mu} A^{\mu}=0 \quad \rightarrow \quad \partial_{\mu} \partial^{\mu} A^{\nu}=j^{\nu}
$$

Coulomb condition

$$
A^{0}=0 \quad \text { or equivalently } \quad \nabla \boldsymbol{A}=0
$$

## The Lagrangian in classical mechanics

In classical mechanics, the Lagrangian is the difference between the kinetic and potential energy: $L(\boldsymbol{q}, \dot{\boldsymbol{q}}) \equiv T-V$. The coordinates $\boldsymbol{q}(t)=\left\{q_{1}(t), \ldots, q_{N}(t)\right\}$ fully describe the system at any given instant $t$. The number $N$ of coordinates is called the number of degrees of freedom of the system.
Let the system move from $A\left(t_{1}\right)$ to $B\left(t_{2}\right)$ along some given path. The action $S$ [path] is defined by the integral of the Lagrangian along the path:

$$
S[\text { path }]=\int_{t_{1}}^{t_{2}} \mathrm{~d} t L(\boldsymbol{q}, \dot{\boldsymbol{q}})
$$

The action $S$ assigns a number to each path and is thus a function of the path. In mathematics, $S$ is called a functional.
The principle of least action states that the system will evolve along the path that minimises the action.
Let $q(t)$ be a path and $q(t)+\delta q(t)$ be some deviating path between the same points $A\left(t_{1}\right)$ and $B\left(t_{2}\right)$. That is, $\delta q\left(t_{1}\right)=\delta q\left(t_{2}\right)=0$. The variation in the action is then given by

$$
\delta S=\int_{t_{1}}^{t_{2}} \mathrm{~d} t \delta L(q, \dot{q})=\int_{t_{1}}^{t_{2}} \mathrm{~d} t\left(\frac{\partial L}{\partial q} \delta q+\frac{\partial L}{\partial \dot{q}} \delta \dot{q}\right)_{\mathrm{Iwant}}^{=} 0 .
$$

Because

$$
\frac{\mathrm{d}}{\mathrm{~d} t}\left(\frac{\partial L}{\partial \dot{q}} \delta q\right)=\left(\frac{\mathrm{d}}{\mathrm{~d} t} \frac{\partial L}{\partial \dot{q}}\right) \delta q+\left(\frac{\partial L}{\partial \dot{q}}\right) \delta \dot{q},
$$

we find, by partial integration,

$$
\delta S=\int_{t_{1}}^{t_{2}} \mathrm{~d} t\left(\frac{\partial L}{\partial q}-\frac{\mathrm{d}}{\mathrm{~d} t} \frac{\partial L}{\partial \dot{q}}\right) \delta q+\int_{t_{1}}^{t_{2}} \mathrm{~d}\left(\frac{\partial L}{\partial \dot{q}} \delta q\right) \underset{\mathrm{Iwant}}{=} 0 .
$$

The second integral vanishes because $\delta q\left(t_{1}\right)=\delta q\left(t_{2}\right)=0$.
The first integral vanishes for all $\delta q$ if and only if the term in brackets vanishes, leading to the Euler-Lagrange equations, for $N$ degrees of freedom:

$$
\frac{\delta S}{\delta q_{i}}=\frac{\mathrm{d}}{\mathrm{~d} t}\left(\frac{\partial L}{\partial \dot{q}_{i}}\right)-\frac{\partial L}{\partial q_{i}}=0 \quad i=1, \ldots, N
$$

Solving the EL equations for a given Lagrangian lead to the equations of motion of the system.

## The Hamiltonian in classical mechanics

If $L$ does not explicitly depend on time we have for the time derivative

$$
\frac{\mathrm{d} L}{\mathrm{~d} t}=\frac{\partial L}{\partial \dot{q}} \ddot{q}+\frac{\partial L}{\partial q} \dot{q}
$$

Substituting $\partial L / \partial q$ from the Euler-Lagrange equations gives

$$
\frac{\mathrm{d} L}{\mathrm{~d} t}=\frac{\partial L}{\partial \dot{q}} \ddot{q}+\frac{\mathrm{d}}{\mathrm{~d} t}\left(\frac{\partial L}{\partial \dot{q}}\right) \dot{q}=\frac{\mathrm{d}}{\mathrm{~d} t}\left(\frac{\partial L}{\partial \dot{q}} \dot{q}\right) \rightarrow \frac{\mathrm{d}}{\mathrm{~d} t}\left(\frac{\partial L}{\partial \dot{q}} \dot{q}-L\right)=0
$$

The term in brackets is the Legendre transform of $L$ and is called the Hamiltonian:

$$
H \stackrel{\text { def }}{=} \frac{\partial L}{\partial \dot{q}} \dot{q}-L=p \dot{q}-L \quad \text { with } \quad p \stackrel{\text { def }}{=} \frac{\partial L}{\partial \dot{q}},
$$

where we have also introduced the canonical momentum $p$. The Hamiltonian is identified with the total energy $E=T+V$ which is thus conserved in the time evolution of the system. This is an example of a conservation law.
In the Lagrangian, the dependence on $\dot{q}$ resides in the kinetic energy term $T$ while the dependence on $q$ is contained in the potential energy $V$. Thus if $V=0$ (or a constant) we have in the EL equations

$$
\frac{\partial L}{\partial q}=0 \quad \rightarrow \quad \frac{\mathrm{~d}}{\mathrm{~d} t}\left(\frac{\partial L}{\partial \dot{q}}\right)=\frac{\mathrm{d} p}{\mathrm{~d} t}=0
$$

Thus the momentum $p$ is conserved in a system that is not under the influence of an external potential. This is another example of a conservation law.

The Hamiltonian equations of motion are

$$
\dot{q}=\frac{\partial H}{\partial p} \quad \text { and } \quad \dot{p}=-\frac{\partial H}{\partial q}
$$

This can be derived as follows. Consider the total differential

$$
\mathrm{d} L=\frac{\partial L}{\partial q} \mathrm{~d} q+\frac{\partial L}{\partial \dot{q}} \mathrm{~d} \dot{q}
$$

Now

$$
\frac{\partial L}{\partial q}=\dot{p}(\text { from EL }), \quad \frac{\partial L}{\partial \dot{q}}=p(\text { by definition }),
$$

and thus, using $p \mathrm{~d} \dot{q}=d(p \dot{q})-\dot{q} \mathrm{~d} p$, we obtain

$$
\mathrm{d} L=\dot{p} \mathrm{~d} q+\mathrm{d}(p \dot{q})-\dot{q} \mathrm{~d} p \quad \rightarrow \quad \mathrm{~d}(p \dot{q}-L)=\mathrm{d} H=\dot{q} \mathrm{~d} p-\dot{p} \mathrm{~d} q,
$$

from which the Hamiltonian equations immediately follow.

## Dirac $\delta$-function

- The Dirac $\delta$-function can be defined by ${ }^{3}$

$$
\delta(x)=\left\{\begin{array}{ll}
0, & \text { if } x \neq 0 \\
\infty, & \text { if } x=0
\end{array} \quad \text { with } \quad \int_{-\infty}^{\infty} \delta(x) \mathrm{d} x=1\right.
$$

- Generalisation to more dimensions is trivial, like $\delta(\boldsymbol{r}) \equiv \delta(x) \delta(y) \delta(z)$.
- For $x \rightarrow 0$ we may write $f(x) \delta(x)=f(0) \delta(x)$ so that

$$
\int_{-\infty}^{\infty} f(x) \delta(x) \mathrm{d} x=f(0) \quad \text { and } \quad \int_{-\infty}^{\infty} f(x) \delta(x-a) \mathrm{d} x=f(a)
$$

- For a linear transformation $y=k(x-a)$ we have

$$
\delta(y)=\frac{1}{|k|} \delta(x-a)
$$

This is straight-forward to prove by showing that $\delta(y)$ satisfies the definition of the $\delta$-function given above.

- Likewise, if $\left\{x_{i}\right\}$ is the set of points for which $f\left(x_{i}\right)=0$, then it is easy to show by Taylor expansion around the $x_{i}$ that

$$
\delta[f(x)]=\sum_{i} \frac{1}{\left|f^{\prime}\left(x_{i}\right)\right|} \delta\left(x-x_{i}\right)
$$

- There exist many representations of the $\delta$-function, for instance,

$$
\delta(\boldsymbol{r})=\frac{1}{(2 \pi)^{3}} \int e^{i \boldsymbol{k} \cdot \boldsymbol{r}} \mathrm{~d}^{3} \boldsymbol{k} \quad \text { or } \quad \delta(x)=\frac{\mathrm{d} \theta(x)}{\mathrm{d} x}
$$

with $\quad \theta(x)=\left\{\begin{array}{l}0, \text { for } x<0 \\ 1, \text { for } x \geq 0\end{array} \quad\right.$ (Heaviside step function).

[^1]
## Green functions

- Let $\Omega$ be some linear differential operator. A Green function of the operator $\Omega$ is a solution of the differential equation

$$
\Omega G(\boldsymbol{r})=\delta(\boldsymbol{r})
$$

These Green functions can be viewed as some potential caused by a point source at $\boldsymbol{r}$.

- Once we have the Green function we can immediately solve the differential equation for any source density $s(\boldsymbol{r})$

$$
\Omega \psi(\boldsymbol{r})=s(\boldsymbol{r})
$$

By substitution it is easy to see that ( $\psi_{0}$ is the solution of $\Omega \psi_{0}=0$ )

$$
\psi(\boldsymbol{r})=\psi_{0}(\boldsymbol{r})+\int G\left(\boldsymbol{r}-\boldsymbol{r}^{\prime}\right) s\left(\boldsymbol{r}^{\prime}\right) \mathrm{d} \boldsymbol{r}^{\prime}
$$

Here it is clearly seen that $G\left(\boldsymbol{r}-\boldsymbol{r}^{\prime}\right)$ 'propagates' the contribution from the source element $s\left(\boldsymbol{r}^{\prime}\right) \mathrm{d} \boldsymbol{r}^{\prime}$ to the potential $\psi(\boldsymbol{r})$.

- A few well-known Green functions are ...

$$
\begin{array}{lrl}
\nabla^{2} G(\boldsymbol{r})=\delta(\boldsymbol{r}) & G(\boldsymbol{r}) & =-1 /(4 \pi r) \\
\left(\nabla^{2}+\boldsymbol{k}^{2}\right) G(\boldsymbol{r})=\delta(\boldsymbol{r}) & G^{ \pm}(\boldsymbol{r}) & =-\exp ( \pm i k r) /(4 \pi r) \\
\left(\nabla^{2}-m^{2}\right) G(\boldsymbol{r})=\delta(\boldsymbol{r}) & G(\boldsymbol{r}) & =-\exp (-m r) /(4 \pi r)
\end{array}
$$

- ... and here are their Fourier transforms

$$
\begin{array}{ll}
G(\boldsymbol{r})=-1 /(4 \pi r) & \tilde{G}(\boldsymbol{q})=-1 / q^{2} \\
G^{+}(\boldsymbol{r})=-\exp (i k r) /(4 \pi r) & \tilde{G}^{+}(\boldsymbol{q})=1 /\left(k^{2}-q^{2}+i \varepsilon\right) \\
G(\boldsymbol{r})=-\exp (-m r) /(4 \pi r) & \tilde{G}(\boldsymbol{q})=-1 /\left(q^{2}+m^{2}\right)
\end{array}
$$

## Non-relativistic scattering theory I

- Classical relation $E=p^{2} / 2 m+V$ with substitution $E \rightarrow i \partial / \partial t$ and $\boldsymbol{p}=-i \nabla$ gives the Schroedinger equation

$$
i \frac{\partial}{\partial t} \psi(\boldsymbol{r}, t)=-\left[\frac{\nabla^{2}}{2 m}-V(\boldsymbol{r}, t)\right] \psi(\boldsymbol{r}, t)
$$

- Separate $\psi(\boldsymbol{r}, t)=\phi(t) \psi(\boldsymbol{r})$. Dividing through by $\phi \psi$ gives

$$
\frac{i \partial_{t} \phi(t)}{\phi(t)}=-\frac{\left[\nabla^{2}-2 m V(\boldsymbol{r}, t)\right] \psi(\boldsymbol{r})}{2 m \psi(\boldsymbol{r})}
$$

Assume now that $V$ does not depend on $t$. The left and right-hand side must then be equal to a constant, say $E$, and we have

$$
\begin{aligned}
\frac{\partial \phi(t)}{\partial t}=-i E \phi(t) & \rightarrow \quad \phi(t)=e^{-i E t} \\
{\left[\nabla^{2}+k^{2}\right] \psi(\boldsymbol{r}) } & =2 m V(\boldsymbol{r}) \psi(\boldsymbol{r})
\end{aligned}
$$

where we have set $k^{2}=2 m E$. Using Green functions we get

$$
\psi(\boldsymbol{r})=\psi_{0}(\boldsymbol{r})-\frac{m}{2 \pi} \int \frac{e^{i k\left|\boldsymbol{r}-\boldsymbol{r}^{\prime}\right|}}{\left|\boldsymbol{r}-\boldsymbol{r}^{\prime}\right|} V\left(\boldsymbol{r}^{\prime}\right) \psi\left(\boldsymbol{r}^{\prime}\right) \mathrm{d} \boldsymbol{r}^{\prime}
$$

- For large $r \gg r^{\prime}$ we have $\left|\boldsymbol{r}-\boldsymbol{r}^{\prime}\right| \approx r-\hat{\boldsymbol{r}} \boldsymbol{r}^{\prime}$ so that

$$
\psi(\boldsymbol{r})=\psi_{0}(\boldsymbol{r})-\frac{m}{2 \pi} \frac{e^{i k r}}{r} \int e^{-i k \hat{r} r^{\prime}} V\left(\boldsymbol{r}^{\prime}\right) \psi\left(\boldsymbol{r}^{\prime}\right) \mathrm{d} \boldsymbol{r}^{\prime}
$$

- We set $\boldsymbol{k}^{\prime} \equiv k \hat{\boldsymbol{r}}$ and write, formally,

$$
\psi(\boldsymbol{r})=\psi_{0}(\boldsymbol{r})+f\left(\boldsymbol{k}^{\prime}\right) \frac{e^{i k r}}{r} \quad \text { with } \quad f\left(\boldsymbol{k}^{\prime}\right) \equiv-\frac{m}{2 \pi} \int e^{-i \boldsymbol{k}^{\prime} \boldsymbol{r}^{\prime}} V\left(\boldsymbol{r}^{\prime}\right) \psi\left(\boldsymbol{r}^{\prime}\right) \mathrm{d} \boldsymbol{r}^{\prime}
$$

The function $f\left(\boldsymbol{k}^{\prime}\right)$ is called the scattering amplitude.

## Non-relativistic scattering theory II

- An incoming plane wave $\psi_{\mathrm{in}}=B e^{i k z}$ describes beam particles moving along the $z$ axis with momentum $k$. The wave function is normalised such that $\rho=\psi^{*} \psi=|B|^{2}$ is the particle density (number of particles per unit volume). The current density is

$$
\boldsymbol{j}_{\text {in }}=\frac{1}{2 m i}\left(\psi^{*} \nabla \psi-\psi \nabla \psi^{*}\right)=|B|^{2} \frac{\boldsymbol{k}}{m}=\rho \frac{\boldsymbol{k}}{m}=\rho \boldsymbol{v}
$$

with $\boldsymbol{v}$ the velocity of the particle. The number of beam particles passing per second through an area $A$ is $R_{\text {in }}=\rho v A=\left|\boldsymbol{j}_{\text {in }}\right| A$. Likewise, the number of scattered particles that pass per second through an area $r^{2} \mathrm{~d} \Omega$ is $R_{\mathrm{sc}}=\left|\boldsymbol{j}_{\mathrm{sc}}\right| r^{2} \mathrm{~d} \Omega$.

- We now imagine a hypothetical area $\mathrm{d} \sigma$ such that the number of beam particles that pass through that area is equal to the number of particles that scatter in the solid angle $\mathrm{d} \Omega$. We then have, by definition, $\left|\boldsymbol{j}_{\text {in }}\right| \mathrm{d} \sigma=\left|\boldsymbol{j}_{\text {sc }}\right| r^{2} \mathrm{~d} \Omega$, or

$$
\frac{\mathrm{d} \sigma}{\mathrm{~d} \Omega}=\frac{r^{2}\left|\boldsymbol{j}_{\mathrm{sc}}\right|}{\left|\boldsymbol{j}_{\text {in }}\right|}
$$

The quantity $\mathrm{d} \sigma / \mathrm{d} \Omega$ is called a differential cross section.

- For our scattered wave $\psi_{\text {sc }}=f\left(\boldsymbol{k}^{\prime}\right) e^{i k r} / r$ we find

$$
j_{\mathrm{sc}}=\frac{1}{2 m i}\left(\psi^{*} \frac{\partial \psi}{\partial r}-\psi \frac{\partial \psi^{*}}{\partial r}\right)=\left|f\left(\boldsymbol{k}^{\prime}\right)\right|^{2} \frac{k}{m r^{2}}
$$

and thus

$$
\frac{\mathrm{d} \sigma}{\mathrm{~d} \Omega}=\left|f\left(\boldsymbol{k}^{\prime}\right)\right|^{2}
$$

where we have assumed $\rho=1$ and $k_{\mathrm{sc}}=k_{\mathrm{in}}$ (elastic scattering).

## Non-relativistic scattering theory III

- Recall that for scattering on a potential, the outgoing wave is

$$
\psi_{\text {out }}(\boldsymbol{r})=\psi_{0}(\boldsymbol{r})+f\left(\boldsymbol{k}^{\prime}\right) \frac{e^{i k r}}{r}
$$

with

$$
f\left(\boldsymbol{k}^{\prime}\right) \equiv-\frac{m}{2 \pi} \int e^{-i \boldsymbol{k}^{\prime} \boldsymbol{r}^{\prime}} V\left(\boldsymbol{r}^{\prime}\right) \psi_{\text {out }}\left(\boldsymbol{r}^{\prime}\right) \mathrm{d} \boldsymbol{r}^{\prime}
$$

Here $\boldsymbol{k}^{\prime}$ is the momentum vector of the scattered particle.

- The problem now is that $\psi_{\text {out }}$ occurs on both sides of the equation above. A first order approximation is achieved by setting in the scattering amplitude $\psi_{\text {out }} \approx \psi_{\text {in }}=e^{i k z}=e^{i k r}$. This gives

$$
f\left(\boldsymbol{k}, \boldsymbol{k}^{\prime}\right) \equiv-\frac{m}{2 \pi} \int e^{i\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right)^{\prime}} V\left(\boldsymbol{r}^{\prime}\right) \mathrm{d} \boldsymbol{r}^{\prime}=-\frac{m}{2 \pi} \int e^{-i \boldsymbol{q} \boldsymbol{r}^{\prime}} V\left(\boldsymbol{r}^{\prime}\right) \mathrm{d} \boldsymbol{r}^{\prime}
$$

where we have set the momentum transfer $\boldsymbol{q} \equiv \boldsymbol{k}^{\prime}-\boldsymbol{k}$. In this so-called Born approximation, the scattering amplitude $f\left(\boldsymbol{k}, \boldsymbol{k}^{\prime}\right) \equiv f(\boldsymbol{q})$ is thus the Fourier transform of the potential.

- Example: Yukawa potential $V(r)=Q_{1} Q_{2} e^{-a r} / r$

$$
\begin{gathered}
f(\boldsymbol{q})=-\frac{m Q_{1} Q_{2}}{2 \pi} \int \frac{e^{-a r^{\prime}}}{r^{\prime}} e^{-i \boldsymbol{q} \boldsymbol{r}^{\prime}} \mathrm{d} \boldsymbol{r}^{\prime}=\cdots=\frac{2 m Q_{1} Q_{2}}{q^{2}+a^{2}} \\
\frac{\mathrm{~d} \sigma}{\mathrm{~d} \Omega}=|f(\boldsymbol{q})|^{2}=\left[\frac{2 m Q_{1} Q_{2}}{q^{2}+a^{2}}\right]^{2}
\end{gathered}
$$

- Example: Coulomb potential $V(r)=Q_{1} Q_{2} / r$ set $a=0$ above:

$$
\frac{\mathrm{d} \sigma}{\mathrm{~d} \Omega}=\left[\frac{2 m Q_{1} Q_{2}}{q^{2}}\right]^{2}=\left[\frac{Q_{1} Q_{2}}{2 m v^{2} \sin ^{2}(\theta / 2)}\right]^{2}
$$

This is the famous formula for Rutherford scattering.

## Dirac's bra-ket notation I

- A state vector $\psi_{\alpha}$ can be represented by a column vector of complex numbers in a Hilbert space and is denoted by the ket $|\alpha\rangle$. To each ket is associated a bra vector $\langle\alpha|$ in a dual Hilbert space. This bra is represented by the conjugate transpose $\psi_{\alpha}^{\dagger}$, that is, by the row vector of complex conjugates. The operation of Hermitian conjugation turns a bra into a ket and vice versa

$$
|\alpha\rangle^{\dagger}=\langle\alpha| \quad \text { and } \quad(c|\alpha\rangle)^{\dagger}=\langle\alpha| c^{*} \quad(c \text { any complex number })
$$

Note that the Hermitian conjugate of a c-number is the complex conjugate. The inproduct $\psi_{\alpha}^{\dagger} \cdot \psi_{\beta}$ is denoted by $\langle\alpha \mid \beta\rangle$ and is a c-number so that

$$
\langle\beta \mid \alpha\rangle \equiv\langle\alpha \mid \beta\rangle^{\dagger}=c^{\dagger}=c^{*}=\langle\alpha \mid \beta\rangle^{*}
$$

- An operator $O$ transforms a ket $|\alpha\rangle$ into another ket, say $|\gamma\rangle$. The operator and its Hermitian conjugate are then defined by

$$
O|\alpha\rangle=|\gamma\rangle \quad \text { and } \quad\langle\alpha| O^{\dagger}=\langle\gamma|
$$

Multiplying from the left with $\langle\beta|$ and from the right with $|\beta\rangle$ we find the relation between the matrix elements of $O$ and $O^{\dagger}$

$$
\begin{aligned}
& O_{\beta \alpha} \equiv\langle\beta| O|\alpha\rangle=\langle\beta \mid \gamma\rangle \\
& O_{\alpha \beta}^{\dagger} \equiv\langle\alpha| O^{\dagger}|\beta\rangle=\langle\gamma \mid \beta\rangle=\langle\beta \mid \gamma\rangle^{*}=\langle\beta| O|\alpha\rangle^{*}=O_{\beta \alpha}^{*}
\end{aligned}
$$

- An operator for which $O=O^{\dagger}$ is called self-adjoint or Hermitian. Observable quantities are always represented by Hermitian operators. Indeed, the expectation value $\langle\alpha| O|\alpha\rangle$ is then real, as it should be, since

$$
\langle\alpha| O|\alpha\rangle \equiv\langle\alpha| O^{\dagger}|\alpha\rangle=\langle\alpha| O|\alpha\rangle^{*}
$$

## Dirac's bra-ket notation II

- An orthonormal basis is written as $\left|e_{i}\right\rangle$ with $\left\langle e_{i} \mid e_{j}\right\rangle=\delta_{i j}$. On this basis, a state $|\alpha\rangle$ is given by the linear combination

$$
|\alpha\rangle=\sum_{i}\left|e_{i}\right\rangle\left\langle e_{i} \mid \alpha\right\rangle
$$

The operator $\left|e_{i}\right\rangle\left\langle e_{i}\right|$ is called a projection operator, for obvious reasons. The closure relation reads $\sum_{i}\left|e_{i}\right\rangle\left\langle e_{i}\right|=1$

- We denote the wave function $\psi_{\alpha}(\boldsymbol{r})$ by $\langle\boldsymbol{r} \mid \alpha\rangle$ and its Hermitian conjugate $\psi_{\alpha}^{\dagger}(\boldsymbol{r})$ by $\langle\alpha \mid \boldsymbol{r}\rangle$. In particular, the wave function of a momentum eigenstate is $\langle\boldsymbol{r} \mid \boldsymbol{k}\rangle \propto e^{i \boldsymbol{k} r}$.
- For the complete set of states $|\boldsymbol{r}\rangle$ the closure relation reads

$$
\int|\boldsymbol{r}\rangle\langle\boldsymbol{r}| \mathrm{d} \boldsymbol{r}=1
$$

From this, we nicely recover the expression for the inproduct of two wave functions

$$
\langle\alpha \mid \beta\rangle=\int\langle\alpha \mid \boldsymbol{r}\rangle\langle\boldsymbol{r} \mid \beta\rangle \mathrm{d} \boldsymbol{r}=\int \psi_{\alpha}^{*}(\boldsymbol{r}) \psi_{\beta}(\boldsymbol{r}) \mathrm{d} \boldsymbol{r}
$$

that of the delta function

$$
\delta\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right)=\left\langle\boldsymbol{k}^{\prime} \mid \boldsymbol{k}\right\rangle=\int\left\langle\boldsymbol{k}^{\prime} \mid \boldsymbol{r}\right\rangle\langle\boldsymbol{r} \mid \boldsymbol{k}\rangle \mathrm{d} \boldsymbol{r} \propto \frac{1}{(2 \pi)^{3}} \int e^{i\left(\boldsymbol{k}-\boldsymbol{k}^{\prime}\right) \boldsymbol{r}} \mathrm{d} \boldsymbol{r}
$$

and also that of Fourier transforms

$$
\psi(\boldsymbol{k})=\langle\boldsymbol{k} \mid \psi\rangle=\int\langle\boldsymbol{k} \mid \boldsymbol{r}\rangle\langle\boldsymbol{r} \mid \psi\rangle \mathrm{d} \boldsymbol{r} \propto \int e^{-i \boldsymbol{k} \boldsymbol{r}} \psi(\boldsymbol{r}) \mathrm{d} \boldsymbol{r}
$$

## Dirac equation

- Dirac equation:

$$
\begin{gathered}
i \gamma^{\mu} \partial_{\mu} \psi-m \psi=0 \\
\underbrace{(\not p-m) u=0}_{\text {particle in }}, \quad \underbrace{\bar{u}(\not p-m)=0}_{\text {particle out }}, \underbrace{(\not p+m) v=0}_{\text {antiparticle out }}, \underbrace{\bar{v}(\not p+m)=0}_{\text {antiparticle in }} \\
\bar{\psi}=\psi^{\dagger} \gamma^{0}, \quad \not x=\gamma^{\mu} a_{\mu}
\end{gathered}
$$

- Pauli matrices:

$$
\begin{gathered}
\sigma_{x}=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right), \quad \sigma_{y}=\left(\begin{array}{rr}
0 & -i \\
i & 0
\end{array}\right), \quad \sigma_{z}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right) \\
\sigma_{i} \sigma_{j}=\delta_{i j}+i \epsilon_{i j k} \sigma_{k}, \quad \sigma_{i}^{\dagger}=\sigma_{i}=\sigma_{i}^{-1}, \quad\left[\sigma_{i}, \sigma_{j}\right]=2 \epsilon_{i j k} \sigma_{k} \\
(\boldsymbol{a} \cdot \boldsymbol{\sigma})(\boldsymbol{b} \cdot \boldsymbol{\sigma})=\boldsymbol{a} \cdot \boldsymbol{b}+i \boldsymbol{\sigma} \cdot(\boldsymbol{a} \times \boldsymbol{b}) \\
\exp (i \boldsymbol{\theta} \cdot \boldsymbol{\sigma})=\cos |\boldsymbol{\theta}|+i(\hat{\boldsymbol{\theta}} \cdot \boldsymbol{\sigma}) \sin |\boldsymbol{\theta}|
\end{gathered}
$$

- Dirac matrices:

$$
\begin{gathered}
\gamma^{0}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right), \quad \gamma^{i}=\left(\begin{array}{cc}
0 & \sigma_{i} \\
-\sigma_{i} & 0
\end{array}\right), \quad \gamma^{5}=i \gamma^{0} \gamma^{1} \gamma^{2} \gamma^{3}=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right) \\
\gamma^{0^{\dagger}}=\gamma^{0}, \quad \gamma^{i \dagger}=-\gamma^{i}, \quad \gamma^{0} \gamma^{\mu \dagger} \gamma^{0}=\gamma^{\mu} \\
\left\{\gamma^{\mu}, \gamma^{\nu}\right\}=2 g^{\mu \nu}, \quad\left\{\gamma^{\mu}, \gamma^{5}\right\}=0, \quad\left(\gamma^{5}\right)^{2}=1
\end{gathered}
$$


[^0]:    ${ }^{2}$ This is the relation between the coordinates $x^{\mu}$ of an event observed in a system $S$ and the coordinates $x^{\prime \mu}$ of that same event observed in a system $S^{\prime}$ that moves with a velocity $+\beta$ along the $x$-axis of $S$.

[^1]:    ${ }^{3} \mathrm{~A}$ more rigorous mathematical definition is usually in terms of a limiting sequence of functions.

