Classical Mechanics

Motion in \mathbb{R}^1

Let x(t) denote position. Then, $v(t) = \frac{dx}{dt} = \dot{x}(t)$ is velocity (where the \cdot denotes derivative with respect to time), $a(t) = \dot{v}(t) = \ddot{x}(t)$, etc.

Considering Newton's second law, $F(x(t)) = m\ddot{x}(t)$, every exact solution requires initial conditions of $x(t_0)$ and $v(t_0)$. Solutions to Newton's second law are known as trajectories.

Considering a spring of constant k, F(x) = -kx yields the differential equation $m\ddot{x} + kx = 0$. The general solution is

$$x(t) = a\cos(\omega t) + b\cos(\omega t),$$

with $\omega = \sqrt{k/m}$ denoting the frequency. The spring is an example of a simple harmonic oscillator.

Conservation of Energy

For a general force function F(x), the kinetic energy is $\frac{1}{2}mv^2$, and the potential energy is

$$V(x) = -\int F(x)dx,$$

meaning $F(x) = -\frac{dV}{dx}$. The total energy is thus found as

$$E(x, v) = \frac{1}{2}mv^2 + V(x).$$

Proposition (Conservation of Energy). If a particle with trajectory x(t) satisfies $m\ddot{x} = F(x)$, then the energy E is conserved.

Proof.

$$\frac{d}{dt}E(x(t),\dot{x}(t)) = \frac{d}{dt}\left(\frac{1}{2}m(\dot{x}(t))^2 + V(x(t))\right)$$
$$= m\dot{x}(t)\ddot{x}(t) + \frac{dV}{dx}\dot{x}(t)$$
$$= \dot{x}(t)\left(m\ddot{x}(t) - F(x(t))\right).$$

By using the conservation of energy, we can reduce the second order differential equation $F(x) = m\ddot{x}$ to a system of first order differential equations in x(t) and v(t) respectively:

$$\frac{dx}{dt} = v(t)$$
$$\frac{dv}{dt} = \frac{1}{m}F(x(t)).$$

If (x(t), v(t)) satisfies this set of equations, then x(t) satisfies Newton's second law. We say the set of all possible (x, v) forms the phase space for the particle in \mathbb{R}^1 .

In phase space, conservation of energy implies that the set of all (x, v) must lie on the level curve of the energy function: $\{(x, v) \mid E(x, v) = E(x_0, v_0)\}.$

Using the conservation of energy, we find that, though Newton's second law is a second order differential equation in time, it is actually a first order differential equation:

$$\frac{m}{2} (\dot{x}(t))^{2} + V(x(t) = E(x(t_{0}), v(t_{0}))$$

$$\dot{x}(t) = \sqrt{\frac{2(E_{0} - V(x(t)))}{m}}$$

Damping

Suppose we also introduce a force that depends on velocity — in the case of a damped simple harmonic oscillator, the equation for force changes from F = -kx to $F = -kx - \gamma \dot{x}$, with $\gamma > 0$. The damping force acts in the opposite direction of velocity, meaning the particle slows down.

The equation of motion is then

$$m\ddot{x} + \gamma \dot{x} + kx = 0.$$

For γ small, the solutions are a sum sines and cosines multiplied by some exponential decay factor, but for γ large, the solutions are only the exponential decay.

Proposition (Energy Conservation in Damped Systems). Suppose a particle moves along x(t) that satisfies $F(x, \dot{x}) = F_1(x) - \gamma \dot{x}$, with $\frac{dV}{dx} = -F_1(x)$ and $\gamma > 0$. Then,

$$\frac{d}{dt}E(x(t),\dot{x}(t)) = -\gamma \dot{x}(t)^{2}.$$

Proof.

$$\frac{d}{dt}E(x(t),\dot{x}(t)) = \dot{x}(t)\left(m\ddot{x}(t) - F_1(x(t))\right)$$
$$= \dot{x}(t)\left(m\ddot{x}(t) - \left(m\ddot{x}(t) + \gamma\dot{x}(t)\right)\right)$$
$$= -\gamma\dot{x}(t)^2$$

Motion in \mathbb{R}^n

The position of a particle $\mathbf{x}=(x_1,\ldots,x_n)$ lends itself to velocity $\mathbf{v}=(v_1,\ldots,v_n)=(\dot{x}_1,\ldots,\dot{x}_n)$, and $\mathbf{a}=(\ddot{x}_1,\ldots,\ddot{x}_n)$. Similar to in \mathbb{R}^1 , Newton's second law is denoted

$$m\ddot{\mathbf{x}} = \mathbf{F}(\mathbf{x}(t), \dot{\mathbf{x}}(t)).$$

Proposition (Conservation of Energy in *n* Dimensions). The energy function

$$E(\mathbf{x}, \dot{\mathbf{x}}) = \frac{1}{2} m \|\dot{\mathbf{x}}\|^2 + V(\mathbf{x})$$

is only satisfied where $\mathbf{F} = -\nabla V$.

Proof.

$$\frac{d}{dt} \left(\frac{1}{2} m \| \dot{\mathbf{x}} \|^2 + V(\mathbf{x}) \right) = m \sum_{j=1}^{n} \dot{x}_j \ddot{x}_j + \sum_{j=1}^{n} \frac{\partial V}{\partial x_j} \dot{x}_j(t)$$

$$= \dot{\mathbf{x}}(t) \left(m \ddot{\mathbf{x}}(t) + \nabla V \right)$$

$$= \dot{x}(t) \left(\mathbf{F}(x) + \nabla V(\mathbf{x}) \right),$$

which is equal to zero only if $-\nabla V = \mathbf{F}$.

If **F** is a smooth \mathbb{R}^n valued function on $U \subset \mathbb{R}^n$, then **F** is conservative if there exists a smooth real-valued function V such that $\mathbf{F} = -\nabla V$.

In other words, \mathbf{F} is conservative if \mathbf{F} is a gradient field, implying that $\nabla \times \mathbf{F} = 0$. If $\mathbf{F}(\mathbf{x}, \mathbf{y}) = -\nabla V(\mathbf{x}) + \mathbf{F}_2(\mathbf{x}, \mathbf{y})$, with $\mathbf{v} \cdot \mathbf{F}_2 = 0$ for all \mathbf{x} and \mathbf{v} , then energy is conserved along a given trajectory.

Systems of Particles

Let $\mathbf{x}^j = \left(x_1^j, x_2^j, \dots, x_n^j\right)$ denote the *j*th particle of a system of *N* particles. Newton's second law is thus reformulated as

$$m_i \ddot{\mathbf{x}}^j = \mathbf{F}^j (\mathbf{x}^1, \dots, \mathbf{x}^N, \dot{\mathbf{x}}^1, \dots, \dot{\mathbf{x}}^N)$$

The total energy is determined by

$$E(\mathbf{x}^1, \dots \mathbf{x}^N, \mathbf{v}^1, \dots, \mathbf{v}^N) = \left(\sum_{j=1}^N \frac{1}{2} m_j \|\mathbf{v}^j\|^2\right) + V(\mathbf{x}^1, \dots, \mathbf{x}^N).$$

Proposition (Conservation of Energy in a System of Particles). The energy function is constant along each trajectory if $\nabla^j V = -\mathbf{F}^j$, where ∇^j denotes the gradient with respect to \mathbf{x}^j .

The force function along a simply connected domain U in \mathbb{R}^{nN} satisfies $\nabla^j V = -\mathbf{F}^j$ if and only if

$$\frac{\partial F_k^j}{\partial x_m^l} = \frac{\partial F_m^l}{\partial x_k^j}$$

for all j, k, l, m.

Proof.

$$\begin{split} \frac{dE}{dt} &= \sum_{j=1}^{N} \left(m_{j} \dot{\mathbf{x}}^{j} \cdot \ddot{\mathbf{x}}^{j} + \nabla^{j} V \cdot \mathbf{x}^{j} \right) \\ &= \sum_{j=1}^{N} \dot{\mathbf{x}}^{j} \left(m_{j} \ddot{\mathbf{x}}^{j} + \nabla^{j} V \right) \\ &= \sum_{j=1}^{N} \dot{\mathbf{x}} \left(\mathbf{F}^{j} + \nabla^{j} V \right), \end{split}$$

which is equal to zero if $\nabla^j V = -\mathbf{F}^j$.

Applying a higher dimension version of $\nabla \times \mathbf{F}$ to each coordinate pair (a, b), we find the identity that shows \mathbf{F} is a gradient field.

Momentum of a System of Particles

The momentum of a particle \mathbf{p}^{j} is defined by

$$\mathbf{p}^j = m_j \dot{\mathbf{x}}^j.$$

Observe that $\frac{d\mathbf{p}^{j}}{dt} = m_{j}\ddot{\mathbf{x}}^{j} = \mathbf{F}^{j}$. The total momentum is then

$$\mathbf{p} = \sum_{i=1}^{N} \mathbf{p}^{j}.$$

Newton's third law, which states "for every action there is an equal and opposite reaction" applies if

•
$$\mathbf{F}^j = \sum_{k \neq j} \mathbf{F}^{j,k}(\mathbf{x}^j, \mathbf{y}^j);$$

$$\bullet \ \mathbf{F}^{j,k}(\mathbf{x}_i,\mathbf{x}_k) = -\mathbf{F}^{k,j}(\mathbf{x}^k,\mathbf{x}^j).$$

If each \mathbf{F}^{j} is also a conservative force, then satisfying these conditions yields potential energy in the form of

$$V(\mathbf{x}^1, \mathbf{x}^2, \dots, \mathbf{x}^N) = \sum_{i < k} V^{j,k} (\mathbf{x}^j - \mathbf{x}^k).$$

Proposition (Newton's Third Law and Conservation of Momentum). *If the system of particles satisfies the conditions of*

$$\bullet \ \mathbf{F}^j = \sum_{k \neq i} \mathbf{F}^{j,k}(\mathbf{x}^j, \mathbf{y}^j)$$

• and
$$\mathbf{F}^{j,k}(\mathbf{x}_j,\mathbf{x}_k) = -\mathbf{F}^{k,j}(\mathbf{x}^k,\mathbf{x}^j)$$
,

then total momentum is conserved.

Proof.

$$\begin{aligned} \frac{d\mathbf{p}}{dt} &= \sum_{j=1}^{N} \mathbf{F}^{j} \\ &= \sum_{j=1}^{N} \sum_{k \neq j} \mathbf{F}^{j,k} (\mathbf{x}^{j}, \mathbf{x}^{k}), \end{aligned}$$

and since $F^{j,k}(\mathbf{x}^j,\mathbf{x}^k) + \mathbf{F}^{k,j}(\mathbf{x}^k,\mathbf{x}^j) = 0$, we find $\frac{d\mathbf{p}}{dt} = 0$.

Proposition (Translation Invariance of Potential). Let V denote the potential for a conservative force. Then, momentum is conserved if and only if V is translation invariant, meaning that for all $\mathbf{a} \in \mathbb{R}^n$,

$$V(\mathbf{x}^1 + \mathbf{a}, \mathbf{x}^2 + \mathbf{a}, \dots, \mathbf{x}^N + \mathbf{a}) = V(\mathbf{x}^1, \mathbf{x}^2, \dots, \mathbf{x}^N)$$

Proof. Let $\mathbf{a} = t\mathbf{e}_k$. Then, differentiating at t = 0 with respect to t, we find

$$0 = \sum_{j=1}^{N} \frac{\partial V}{\partial x_k^j}$$
$$= -\sum_{j=1}^{N} F_k^j$$
$$= -\sum_{j=1}^{N} \frac{dp_k^j}{dt}$$
$$= -\frac{dp_k}{dt},$$

with p_k denoting the kth component of **p**. Therefore, **p** is constant in time.

If \mathbf{p} is conserved, then the sum of all forces is 0 at each point for all t, meaning that for all t,

$$\frac{d}{dt}V(\mathbf{x}^1 + t\mathbf{a}, \mathbf{x}^2 + t\mathbf{a}, \dots, \mathbf{x}^N + t\mathbf{a}) = \sum_{j=1}^N \nabla^j V(\mathbf{x}^1 + t\mathbf{a}, \mathbf{x}^2 + t\mathbf{a}, \dots, \mathbf{x}^n + t\mathbf{a}) \cdot \mathbf{a}$$

$$= -\left(\sum_{j=1}^N \mathbf{F}^j(\mathbf{x}^1 + t\mathbf{a}, \mathbf{x}^2 + t\mathbf{a}, \dots, \mathbf{x}^N + t\mathbf{a})\right) \cdot \mathbf{a}$$

$$= 0$$

meaning V is equal at t = 0 and t = 1.

Center of Mass

For a system of N particles, the center of mass is denoted

$$\mathbf{c} = \sum_{j=1}^{N} \frac{m_j}{\sum_{j=1}^{N} m_j} \mathbf{x}_j.$$

We denote $\sum_{j=1}^{N} m_j = M$. Differentiating **c**, we get

$$\frac{d\mathbf{c}}{dt} = \frac{1}{M} \sum_{j=1}^{N} m_j \dot{\mathbf{x}}^j$$
$$= \frac{\mathbf{p}}{M}.$$

Notice that if **p** is conserved, then $\mathbf{c}(t) = \mathbf{c}(t_0) + (t - t_0) \frac{\mathbf{p}}{M}$.

For a system of two particles, if $V(\mathbf{x}^1, \mathbf{x}^2)$ is invariant under translation, then $V(\mathbf{x}^1, \mathbf{x}^2) = \tilde{V}(\mathbf{x}^1 - \mathbf{x}^2)$, and $\tilde{V}(\mathbf{a}) = V(\mathbf{a}, 0)$.

The positions \mathbf{x}^1 and \mathbf{x}^2 can be recovered from knowledge about \mathbf{c} and the relative position $\mathbf{y} := \mathbf{x}^1 - \mathbf{x}^2$:

$$\mathbf{x}^1 = \frac{\mathbf{c} + m_2 \mathbf{y}}{m_1 + m_2}$$
$$\mathbf{x}^2 = \frac{\mathbf{c} - m_1 \mathbf{y}}{m_1 + m_2}.$$

Thus, we can calculate

$$\begin{split} \ddot{\mathbf{y}} &= \ddot{\mathbf{x}}^1 - \ddot{\mathbf{x}}^2 \\ &= -\frac{1}{m_1} \nabla \tilde{V} \left(\mathbf{x}^1 - \mathbf{x}^2 \right) - \frac{1}{m_2} \nabla \tilde{V} \left(\mathbf{x}^1 - \mathbf{x}^2 \right). \end{split}$$

Motion of Relative Position under Translation Invariant Potential

For a two particle system with translation invariant potential, the relative position $\mathbf{y} = \mathbf{x}^1 - \mathbf{x}^2$ is a solution to the differential equation

$$\mu\ddot{\mathbf{y}} = -\nabla \tilde{V}(\mathbf{y})$$
,

where

$$\mu = \frac{m_1 m_2}{m_1 + m_2}.$$

This implies that when momentum is conserved, the relative position of the two particle system evolves as a one-particle system with effective mass μ .

Angular Momentum

A particle moving in \mathbb{R}^2 with position \mathbf{x} , velocity \mathbf{v} , and momentum $\mathbf{p} = m\mathbf{v}$ has angular momentum J denoted as

$$J = x_1 p_2 - x_2 p_1,$$

or $J = \|\mathbf{x} \times \mathbf{p}\| = \|\mathbf{x}\| \|\mathbf{p}\| \sin \phi$, with ϕ measured counterclockwise. In polar coordinates, we find

$$J = mr^2 \frac{d\theta}{dt}$$
$$= 2M \frac{dA}{dt},$$

where $A = 1/2 \int r^2 d\theta$ denotes the area swept out by $\mathbf{x}(t)$.

Proposition (Conservation of Angular Momentum). Suppose a particle of mass m is moving in \mathbb{R}^2 under the influence of a conservative force with potential $V(\mathbf{x})$. V is invariant under rotation if and only if J is conserved.

Proof.

$$\frac{dJ}{dt} = \frac{dx_1}{dt}p_2 + x_1\frac{dp_2}{dt} - \frac{dx_2}{dt}p_1 - x_2\frac{dp_1}{dt}$$

$$= \frac{1}{m}p_1p_2 - x_1\frac{\partial V}{\partial x_2} - \frac{1}{m}p_2p_1 + x_2\frac{\partial V}{\partial x_1}$$

$$= x_2\frac{\partial V}{\partial x_1} - x_1\frac{\partial V}{\partial x_2}.$$

Alternatively, consider $R_{\theta} = \begin{bmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{bmatrix}$. Differentiating V along R_{θ} , we get

$$\frac{d}{d\theta}V(R_{\theta}\mathbf{x})\Big|_{\theta=0} = \frac{\partial V}{\partial x}\frac{dx}{d\theta} + \frac{\partial V}{\partial y}\frac{dy}{d\theta}$$
$$= -x_2\frac{\partial V}{\partial x_1} + x_1\frac{\partial V}{\partial x_2}$$
$$= -\frac{dJ}{dt}(\mathbf{x})$$

Thus, $\frac{dJ}{dt} = 0$ if and only if the angular derivative of V is zero.

As a result of the conservation of angular momentum, we thus get Kepler's Second Law: if $\mathbf{x}(t)$ is the trajectory of a particle under the influence of a force with rotationally invariant potential, then the area swept out by $\mathbf{x}(t)$ between t=a and t=b is $\frac{b-a}{2m}J$.

In \mathbb{R}^3 , **J** is a vector given by $\mathbf{x} \times \mathbf{p}$. Meanwhile, in \mathbb{R}^n , the angular momentum is a skew-symmetric matrix defined by

$$J_{jk} = x_j p_k - x_k p_j.$$

The total angular momentum of a system of N particles in \mathbb{R}^n is given by **J** with entries

$$J_{jk} = \sum_{l=1}^{N} (x_{j}^{l} p_{k}^{l} - x_{k}^{l} - p_{j}^{l}).$$

Similar to the case of linear momentum, angular momentum is constant in the presence of a conservative force if and only if the potential function V is rotationally invariant. That is,

$$V(R\mathbf{x}^1, R\mathbf{x}^2, \dots, R\mathbf{x}^N) = V(\mathbf{x}^1, \mathbf{x}^2, \dots, \mathbf{x}^N)$$

for all rotation matrices R.

Hamiltonian Mechanics

The Hamiltonian is the total energy function, but formulated in terms of position and momentum rather than position and velocity. If a particle in \mathbb{R}^n has the usual energy function, we write

$$H(\mathbf{x},\mathbf{p}) = \frac{1}{2m} \sum_{i=1}^{n} p_j^2 + V(\mathbf{x}),$$

where $p_i = m_i \dot{x}_i$. Observe that the equations of motion can be written as

$$\frac{dx_j}{dt} = \frac{\partial H}{\partial p_j}$$
$$\frac{dp_j}{dt} = -\frac{\partial H}{\partial x_j}.$$

In the basic formulation, we can see that the first equation is just $\dot{x}_j = p_j/m$, and $\dot{p}_j = F_j$. The equations of motion written with Hamiltonians are known as Hamilton's equations.

Poisson Bracket

Let f and g be two smooth functions on \mathbb{R}^{2n} , with each element of \mathbb{R}^{2n} being denoted by (\mathbf{x}, \mathbf{p}) . The Poisson bracket of f and g is equal to

$$\{f,g\}\left(\mathbf{x},\mathbf{p}\right) = \sum_{i=1}^{n} \left(\frac{\partial f}{\partial x_{j}} \frac{\partial g}{\partial p_{j}} - \frac{\partial f}{\partial p_{j}} \frac{\partial g}{\partial x_{j}}\right).$$

The Poisson bracket satisfies the following properties:

- Linearity: $\{f, g + ch\} = \{f, g\} + c\{f, h\}$
- Antisymmetry: $\{g, f\} = -\{f, g\}$
- Product Rule: $\{f, gh\} = \{f, g\} h + g \{f, h\}$
- Jacobi Identity: $\{f, \{g, h\}\} + \{h, \{f, g\}\} + \{g, \{h, f\}\} = 0$.

It can be easily verified that the following Poisson bracket relations hold:

$$\{x_j, x_k\} = 0$$

$$\{p_j, p_k\} = 0$$

$$\{x_i, p_k\} = \delta_{ik},$$

where δ_{jk} denotes the Kronecker delta function.

Proposition (Solutions of Hamilton's Equations). If $(\mathbf{x}(t), \mathbf{p}(t))$ is a solution to Hamilton's Equations, then for any smooth f on \mathbb{R}^{2n} , we have

$$\frac{df}{dt} = \{f, h\}.$$

Proof.

$$\frac{df}{dt} = \sum_{j=1}^{n} \left(\frac{\partial f}{\partial x_j} \frac{dx_j}{dt} + \frac{\partial f}{\partial p_j} \frac{dp_j}{dt} \right)$$
$$= \sum_{j=1}^{n} \left(\frac{\partial f}{\partial x_j} \frac{\partial H}{\partial p_j} + \frac{\partial f}{\partial p_j} \left(-\frac{\partial H}{\partial x_j} \right) \right)$$
$$= \{f, H\}.$$

Conserved Quantities

Let $f \in C^1(\mathbb{R}^{2n})$ be called conserved if $f(\mathbf{x}(t), \mathbf{p}(t))$ is independent of t for each solution to Hamilton's equation. Then, f is a conserved quantity if and only if

$${f, H} = 0.$$

Note that H is also a conserved quantity.

Flow and Liouville's Theorem

Solving Hamilton's equations on \mathbb{R}^{2n} yields a flow Φ_t^{I} with $\Phi_t(\mathbf{x}, \mathbf{p})$ equal to the solution at time t with initial condition (\mathbf{x}, \mathbf{p}) .

The Φ_t aren't necessarily defined on all of \mathbb{R}^{2n} , but if Φ_t is defined on \mathbb{R}^{2n} for all t, then we say Φ_t is complete.

Proposition (Liouville^{II}). The Hamiltonian flow preserves the 2n-dimensional measure.

$$dx_1 dx_2 \cdots dx_n dp_1 dp_2 \cdots dp_n$$
.

More specifically, if E is a measurable subset of the domain of Φ_t , then $\mu(\Phi_t(E)) = \mu(E)$.

Proof. Hamilton's equations can be written as

$$\frac{d}{dt} \begin{vmatrix} x_1 \\ \vdots \\ x_n \\ p_1 \\ \vdots \\ p_n \end{vmatrix} = \begin{bmatrix} \frac{\partial H}{\partial p_1} \\ \vdots \\ \frac{\partial H}{\partial p_n} \\ -\frac{\partial H}{\partial x_1} \\ \vdots \\ -\frac{\partial H}{\partial x_n} \end{bmatrix}.$$

Hamilton's equations describe the flow along the vector field appearing on the right side — by a result in vector calculus, III the flow preserves the 2n-dimensional area measure if and only if the divergence of the vector field is zero.

$$\nabla \cdot \begin{bmatrix} \frac{\partial H}{\partial p_1} \\ \vdots \frac{\partial H}{\partial p_n} \\ -\frac{\partial H}{\partial x_1} \\ \vdots \\ -\frac{\partial H}{\partial x_n} \end{bmatrix} = \sum_{k=1}^n \frac{\partial}{\partial x_k} \frac{\partial H}{\partial p_k} - \frac{\partial}{\partial p_k} \frac{\partial H}{\partial x_k}$$
$$= \sum_{k=1}^n \frac{\partial^2 H}{\partial x_k \partial p_k} - \frac{\partial^2 H}{\partial p_k \partial x_k}$$
$$= 0$$

The condition of zero divergence is equivalent to Φ_t preserving a particular symplectic form ω defined by

$$\omega((\mathbf{x}, \mathbf{p}), (\mathbf{x}', \mathbf{p}')) = \mathbf{x} \cdot p' - \mathbf{p} \cdot x',$$

Ithe Φ_t are diffeomorphisms, or differentiable isomorphisms with differentiable inverses

IInot from complex analysis

III Author's Note: I do not know this result yet, but hopefully I will soon!

meaning that for any t and any $(\mathbf{x}, \mathbf{p}) \in \mathbb{R}^{2n}$, the partial derivatives of Φ_t preserves ω .

Alternatively, this is equivalent to Φ_t preserving Poisson brackets:

$$\{f \circ \Phi_t, g \circ \Phi_t\} = \{f, g\} \circ \Phi_t.$$

Thus, Φ_t is an example of a symplectomorphism.

Hamiltonian Flow and Hamiltonian Generators

We say $f \in C^1(\mathbb{R}^{2n})$ is the Hamiltonian generator of the flow that results from solving Hamilton's equations with f substituted for H:

$$\frac{dx_j}{dt} = \frac{\partial f}{\partial p_j}$$
$$\frac{dp_j}{dt} = -\frac{\partial f}{\partial x_j}.$$

It is possible to see that

$$f_{a}(\mathbf{x}, \mathbf{p}) = \mathbf{a} \cdot \mathbf{p}$$

yields the flow

$$\mathbf{x}(t) = \mathbf{x}_0 + t\mathbf{a}$$
$$\mathbf{p}(t) = \mathbf{p}_0,$$

and

$$g_{\rm b}(\mathbf{x},\mathbf{p}) = \mathbf{b} \cdot \mathbf{x}$$

yields the flow

$$\mathbf{x}(t) = \mathbf{x}_0$$
$$\mathbf{p}(t) = \mathbf{p}_0 - t\mathbf{b}.$$

Thus, the Hamiltonian flow generated by momentum yields translation in position, and the Hamiltonian flow generated by position yields translation in momentum.

In this light, we can think of *the* Hamiltonian as the Hamiltonian generator that yields time evolution. Other Hamiltonian generators represent some other family of symmetries of the system.

Proposition (Hamiltonian Flow generated by Angular Momentum). For a particle moving in \mathbb{R}^2 , the Hamiltonian flow generated by

$$J(\mathbf{x}, \mathbf{p}) = x_1 p_2 - x_2 p_1$$

consists of simultaneous rotations of \mathbf{x} and \mathbf{p} .

$$\begin{bmatrix} x_1(t) \\ x_2(t) \end{bmatrix} = \begin{bmatrix} \cos t & -\sin t \\ \sin t & \cos t \end{bmatrix} \begin{bmatrix} x_1(0) \\ x_2(0) \end{bmatrix}$$
$$\begin{bmatrix} p_1(t) \\ p_2(t) \end{bmatrix} = \begin{bmatrix} \cos t & -\sin t \\ \sin t & \cos t \end{bmatrix} \begin{bmatrix} p_1(0) \\ p_2(0) \end{bmatrix}.$$

Proof. Plugging J Hamilton's equations, we get

$$\frac{dx_1}{dt} = \frac{\partial J}{\partial p_1} = -x_2$$

$$\frac{dp_1}{dt} = -\frac{\partial J}{\partial x_1} = -p_2$$

$$\frac{dx_2}{dt} = \frac{\partial J}{\partial p_2} = x_1$$

$$\frac{dp_2}{dt} = -\frac{\partial J}{\partial x_2} = p_1.$$

It's important to note that the parameter t in the Hamiltonian flow for J is the rotation, not time. That is, J is the Hamiltonian generator of rotations.

If f is any smooth function, it is the case that the time derivative of any other function g along the Hamiltonian flow generated by f is $\frac{dg}{dt} = \{g, f\}$. In particular, the derivative of H along the flow generated by f is $\{H, f\}$, meaning that f is constant along the flow generated by H if and only if $\{f, H\} = 0$, which is true if and only if H is constant along the flow generated by H.

Thus, we find that f is conserved for solutions of Hamilton's equations if and only if H is invariant under the Hamiltonian flow generated by f. Of particular note, we find that J is conserved if and only if H is invariant under rotations of \mathbf{x} and \mathbf{p} .

Introduction to Quantum Mechanics

Observable quantities such as position and momentum in quantum mechanics are represented by operators on a complex-valued Hilbert space (an inner product space that is complete with respect to the induced metric) — specifically, these quantities are *unbounded* linear operators.

In physics, the inner product is linear in the second factor and conjugate linear in the first factor:

$$\langle \phi, \lambda \psi \rangle = \lambda \langle \phi, \psi \rangle$$

 $\langle \lambda \phi, \psi \rangle = \overline{\lambda} \langle \phi, \psi \rangle$.

Alternatively, in Dirac notation:

$$\langle \phi \mid \lambda \psi \rangle = \lambda \langle \phi \mid \psi \rangle$$
$$\langle \lambda \phi \mid \psi \rangle = \overline{\lambda} \langle \phi \mid \psi \rangle.$$

A Taste of Operator Theory

A linear operator $A: \mathbf{H} \to \mathbf{H}$ is bounded if it has finite operator norm:

$$\sup_{\|\psi\| \le 1} \|A\psi\| < \infty.$$

For each bounded operator A, there exists a unique bounded operator A^* such that $\langle \phi, A\psi \rangle = \langle A^*\phi, \psi \rangle$. The existence of A^* follows from the Riesz representation theorem.

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^{IV}There is another section on Kepler's Laws in the chapter on Classical Mechanics that I didn't really read in depth. I might include it in the future.

 $^{^{}m V}$ I'm using more operator-theoretic language than the book uses because I'm <u>pretentions</u> a mathematician, not a physicist.

A bounded operator is said to be self-adjoint if $A^* = A$. Self-adjoint operators are nice for a variety of reasons, and as a result we desire for our operators in quantum mechanics to be self-adjoint. However, this brings a significant problem — unbounded self-adjoint operators are not necessarily defined on \mathbf{H} .

We define unbounded operators as linear operators defined on a dense subspace of H:

$$A: \mathsf{Dom}(A) \subseteq \mathbf{H} \to \mathbf{H}$$

subject to

$$\overline{\mathsf{Dom}(A)} = \mathbf{H}.$$

In addition to the domain of A not necessarily being equal to \mathbf{H} , the linear functional $\langle \phi, A \cdot \rangle$ is not necessarily bounded (meaning we cannot use the Riesz representation theorem to find $A^*\phi$). The adjoint of A, as a result, will be defined on a subspace of \mathbf{H} .

A vector $\phi \in \mathbf{H}$ is said to belong to the domain $\mathsf{Dom}(A^*)$ if the linear functional $\langle \phi, A \cdot \rangle$ on $\mathsf{Dom}(A)$ is bounded. Then, we define A^* to be the unique vector χ such that $\langle \chi, \psi \rangle = \langle \phi, A\psi \rangle$ for all $\psi \in \mathsf{Dom}(A)$.

Having defined adjoints of an unbounded operator, we can now commit to defining self-adjoint operators. The operator A is symmetric if $\langle \phi, A\psi \rangle = \langle A\phi, \psi \rangle$ — a symmetric operator is self-adjoint if $\mathsf{Dom}(A) = \mathsf{Dom}(A^*)$ and $A^*\phi = A\phi$ for all $\phi \in \mathsf{Dom}(A)$. Finally, A is essentially self-adjoint if the closure of the graph of A in $\mathbf{H} \times \mathbf{H}$ is self-adjoint.

In sum, A is self-adjoint if A and A^* are the same operator with the same domain, more or less.

Definition (Properties of Symmetric Operators). Let A be a symmetric operator on \mathbf{H} . Then, the following hold:

- (1) For all $\psi \in \text{Dom}(A)$, the quantity $\langle \psi, A\psi \rangle$ is real. More generally, if $\psi, A\psi, ..., A^{m-1}\psi$ belong to Dom(A), then $\langle \psi, A^m \psi \rangle$ is real.
- (2) Suppose λ is an eigenvector for A. Then, $\lambda \in \mathbb{R}$.

Proof. (1) Since A is symmetric,

$$\langle \psi, A\psi \rangle = \langle A\psi, \psi \rangle$$

= $\overline{\langle \psi, A\psi \rangle}$,

for all $\psi \in \text{Dom}(A)$. Similarly, if ψ , $A\psi$, ..., $A^{m-1}\psi \in \text{Dom}(A)$, then we use the symmetry of A to show that

$$\langle \psi, A^m \psi \rangle = \langle A^m \psi, \psi \rangle$$
$$= \overline{\langle \psi, A^m \psi \rangle}.$$

(2) If ψ is an eigenvector for A with eigenvalue λ , then

$$\lambda \langle \psi, \psi \rangle = \langle \psi, A\psi \rangle$$
$$= \langle A\psi, \psi \rangle$$
$$= \overline{\lambda} \langle \psi, \psi \rangle.$$

Since ψ is nonzero by definition, it must be the case that $\lambda = \overline{\lambda}$.

In physical terms, $\langle \psi, A\psi \rangle$ represents the expected value for measurements of A in the state ψ , with λ representing a possible value of this measurement. This is why we want both numbers to be real.

A self-adjoint A allows us to use the spectral theorem to assign each $\psi \in \mathbf{H}$ a probability measure on the real numbers.

Position and Momentum Operators

Consider a particle moving along the real line with wave function $\psi: \mathbb{R} \to \mathbb{C}$. Although ψ will evolve over time, let the particle be fixed in time for now.

We want to define ψ to be a unit vector in $L^2(\mathbb{R})$, meaning

$$\int_{\mathbb{D}} |\psi(x)|^2 dx = 1.$$

The probability that the position of the particle belongs to some $E \subseteq \mathbb{R}$ is

$$\int_{F} |\psi(x)|^2 dx,$$

where E is necessarily a Lebesgue-measurable set.

The expectation value of the position is thus

$$E(x) = \int_{\mathbb{R}} x |\psi(x)|^2 dx,$$

assuming the convergence of the integral, and the mth moment of the position is calculated as

$$E(x^m) = \int_{\mathbb{R}} x^m |\psi(x)|^2 dx,$$

again assuming convergence of the integral.

Definition (Position Operator). The position operator is defined as $X = M_x$, meaning $(X\psi)(x) = x\psi(x)$. With this in mind, we can then see that

$$E(x) = \langle \psi, X\psi \rangle$$

under the standard inner product on $L^2(\mathbb{R})$. The expectation value of X for the state ψ is denoted $\langle X \rangle_{\psi} := \langle \psi, X \psi \rangle$.

The higher moments of position are similarly defined:

$$E(x^m) = \langle \psi, X^m \psi \rangle$$

where X^m denotes *m*-degree composition of X.

Since X is an unbounded linear operator, it is not necessarily the case that $X\psi \in L^2(\mathbb{R})$ if $\psi \in L^2(\mathbb{R})$.

Momentum is encoded in the oscillations of the wave function — the de Broglie hypothesis provides a special relationship between the frequency of oscillation (as a function of position at a fixed time) and the momentum.

Proposition (De Broglie Hypothesis). If the wave function of a particle has angular frequency k, then the momentum p of the particle is

$$p = \hbar k$$

where ħ denotes Planck's constant.

VII don't like this notation either.

VIIOther famous examples of unbounded linear operators on Hilbert spaces include the derivative operator on $A^2(\mathbb{D})$, the space of holomorphic functions on the complex unit disc. I'm doing research on properties of variations of this space.

To be more precise, the de Broglie hypothesis applies to wave functions of the form $\psi(x) = e^{ikx}$, which represents particles that have momentum $p = \hbar k$. Let's develop this a bit further.

Since e^{ikx} is not square integrable over \mathbb{R} , we instead move to the circle, where ψ has period 2π over \mathbb{R} , and

$$\int_0^{2\pi} |\psi(x)|^2 = 1.$$

For any integer k, the normalized wave function $\frac{e^{ikx}}{\sqrt{2\pi}}$ will represent the particle with momentum $p=\hbar k$. This momentum value is definite; that is, $p=\hbar k$ with probability 1 for a particle with wave function $\frac{e^{ikx}}{\sqrt{2\pi}}$.

Of note, the functions $\left\{\frac{e^{ikx}}{\sqrt{2\pi}}\right\}$ for $k \in \mathbb{Z}$ form an orthonormal basis for the Hilbert space of square integrable functions with period 2π .

The wave functions for particles on a circle are thus all of the form

$$\psi(x) = \sum_{k=-\infty}^{\infty} a_k \frac{e^{ikx}}{\sqrt{2\pi}},$$

where the sum is convergent in $L^2([0, 2\pi])$. VIII If ψ is a unit vector, then IX

$$\|\psi\|_{L^2([0,2\pi])}^2 = \sum_{k=-\infty}^{\infty} |a_k|^2$$
= 1

For a particle with wave function ψ expressed as a Fourier series, the momentum isn't definite. We will have to consider that measurement will yield one of the values of $\hbar k$ with probability $|a_k|^2$.

$$E(p) = \sum_{k=-\infty}^{\infty} \hbar k |a_k|^2,$$

with higher moments defined by

$$E(p^m) = \sum_{k=-\infty}^{\infty} (\hbar k)^m |a_k|^2,$$

assuming absolute convergence.

Definition (Momentum Operator). To encode P, our momentum operator, such that for the wave function $\psi \in L^2([0, 2\pi])$, $E(p^m) = \langle \psi, P^m \psi \rangle$, we must have P satisfying

$$Pe^{ikx} = \hbar ke^{ikx}$$

Thus, we would assume that

$$P = -i\hbar \frac{d}{dx}.$$

Even on the real line, we would still expect $P=-i\hbar\frac{d}{dx}$; though e^{ikx} is not square integrable, we can represent any wave function $\psi\in L^2(\mathbb{R})$ as an integral with the Fourier transform:

$$\psi(x) = \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{ikx} \hat{\psi}(k) dk,$$
$$\hat{\psi}(x) = \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-ikx} \psi(x) dx.$$

VIII You may recognize these as Fourier series.

^{IX}We use Parseval's identity to relate the L^2 norm of ψ to the ℓ^2 norm of $\{a_k\}_{k\in\mathbb{Z}}$. Try proving it yourself! Hint: use the Pythagorean theorem.

Plancherel's theorem^X states that the Fourier transform is unitary — that is,

$$\int_{-\infty}^{\infty} |\psi(x)|^2 dx = \int_{-\infty}^{\infty} \left| \hat{\psi}(k) \right|^2 dk = 1.$$

We can imagine $\hat{\psi}(k)$ as the probability density for the momentum of the particle^{XI}.

Thus, we have defined the momentum operator as

$$P = -i\hbar \frac{d}{dx},$$

with, for sufficiently nice $\psi \in L^2(\mathbb{R})$

$$E(p^{m}) = \langle \psi, P^{m} \psi \rangle$$
$$= \int_{-\infty}^{\infty} (\hbar k)^{m} |\hat{\psi}(k)|^{2} dk$$

for all positive integers m.

Proposition (Commutator of Position and Momentum).

$$[X, P] = i\hbar I$$
,

a relation known as the canonical commutation relation.

Proof.

$$PX\psi = -i\hbar \frac{d}{dx} (x\psi(x))$$
$$= -i\hbar \psi(x) - i\hbar x \frac{d\psi}{dx}$$
$$= -i\hbar l\psi + XP\psi$$

Remark. Note the parallel between $\{x, p\} = 1$ in the Poisson bracket and $[X, P] = i\hbar I$ in the commutator.

Proposition (Symmetry of Position and Momentum Operators). For all sufficiently nice ϕ and ψ in $L^2(\mathbb{R})$,

$$\langle \phi, X\psi \rangle = \langle X\phi, \psi \rangle$$

 $\langle \phi, P\psi \rangle = \langle P\phi, \psi \rangle$.

Proof. Let $\phi, \psi \in L^2(\mathbb{R})$ with $x\psi(x), x\phi(x) \in L^2(\mathbb{R})$. Then, since $x \in \mathbb{R}$,

$$\int_{-\infty}^{\infty} \overline{\phi(x)} x \psi(x) dx = \int_{-\infty}^{\infty} \overline{x \phi(x)} \psi(x) dx,$$

with both integrals convergent.

Meanwhile, assume ϕ and ψ are continuously differentiable, vanish at $\pm \infty$, and ϕ , ψ , $\frac{d\phi}{dx}$, $\frac{d\psi}{dx} \in L^2(\mathbb{R})$. Note that $\frac{d\overline{\phi}}{dx} = \frac{\overline{d\phi}}{\overline{dx}}$. Then, integrating by parts,

$$-i\hbar \int_{-n}^{n} \overline{\phi(x)} \frac{d\psi}{dx} dx = -i\hbar \overline{\phi(x)} \psi(x) \Big|_{-n}^{n} + i\hbar \int_{-n}^{n} \overline{\frac{d\phi}{dx}} \psi(x) dx,$$

 $^{^{\}rm X}$ also known as Parseval's theorem when applied to Fourier series rather than the Fourier transform

 $^{^{\}rm XI}$ well, for p/\hbar , but that's basically the same

meaning

$$\begin{split} \int_{-\infty}^{\infty} \overline{\phi(x)} \left(-i\hbar \frac{d\psi}{dx} \right) dx &= i\hbar \int_{-\infty}^{\infty} \overline{\frac{d\phi}{dx}} \psi(x) dx \\ &= \int_{-\infty}^{\infty} \overline{\left(-i\hbar \frac{d\phi}{dx} \right)} \psi(x) dx. \end{split}$$

Thus, we have shown that X and P are symmetric operators on certain dense subspaces of $L^2(\mathbb{R})$. We will have to wait until later to prove that X and P are essentially self-adjoint.

Kinematic Axioms of Quantum Mechanics

These aren't really axioms^{XII} (as in, first principles from which quantum mechanics is derived), but they're essential principles of quantum mechanics. For ease of use, I'm labeling them with my personal interpretation of their content (rather than the numbering system used in the book).

Axiom (Principle of Representation). The state of a quantum mechanical system is represented by $\psi \in \mathbf{H}$ for some Hilbert space \mathbf{H} . If ψ_1 and ψ_2 are two unit vectors in \mathbf{H} with $\psi_1 = c\psi_2$ for some $c \in \mathbb{C}$, ψ_1 and ψ_2 represent the same state.

Each ψ represents a pure state — there are mixed states, but those will be discussed later.

Axiom (Principle of Correspondence). To each real-valued f on the classical phase space, there is an unbounded, self-adjoint operator \hat{f} on the Hilbert space.

For a particle in \mathbb{R}^1 , the phase space in \mathbb{R}^2 is represented by (x,p) for x position and p momentum; the analogue to the classical phase space in quantum mechanics is $L^2(\mathbb{R})$ with the position function f(x,p)=x being M_x (the position operator), and the momentum function g(x,p)=p being $P=-i\hbar\frac{d}{dx}$ (the momentum operator).

Axiom (Principle of Measurement). If a system is in a state represented by $\psi \in \mathbf{H}$, then the measurement for the mth moment of f satisfies

$$E(f^m) = \langle \psi, (\hat{f})^m \psi \rangle.$$

In particular, the measurement for f is

$$E(f) = \langle \psi, \hat{f} \psi \rangle$$
.

Remark. Note that in the quantum system, we are measuring the classical f. However, rather than a definitive value, we need to find the expectation of the operator \hat{f} using the probabilities derived from ψ .

Since \hat{f} is self-adjoint, $E(f^m)$ is real, we can construct a probability measure $\mu_{A,\psi}$ that is the probability distribution for measurements of A in the state ψ .

Proposition (Eigenvectors). If a quantum system is in a state described by $\psi \in \mathbf{H}$, and $\hat{f}\psi = \lambda \psi$ for some $\lambda \in \mathbb{R}^{|X|||}$, then

$$E(f^m) = \lambda^m$$
.

The aforementioned probability measure consistent with this condition is δ_{λ} (where f has value λ with probability 1).

XIIPhysicists once again caught stealing the valor of mathematicians.

 $^{^{\}mbox{\scriptsize XIII}}\mbox{Recall}$ that the eigenvalues of self-adjoint operators are real.

Essentially, if ψ is an eigenvector for \hat{f} , then measurements for f are determined. Thus, we would need to find the probability measure such that

$$\int_{\mathbb{R}} x^m d\mu = \lambda^m,$$

which only works if $\mu = \delta_{\lambda}$.

Note that if ψ is a linear combination of eigenvectors for \hat{f} , then the measurements of f are not deterministic. XIV

Example. Suppose \hat{f} has an orthonormal basis $\{e_j\}$ with distinct real eigenvalues λ_j . Let ψ be a unit vector in \mathbf{H} with the expansion

$$\psi = \sum_{i=1}^{\infty} a_i e_i.$$

Then, the measurement in the state ψ of f will necessarily be some value of λ_j with probability $|a_j|^2$.

Axiom (Principle of Wave Function Collapse). Suppose a quantum system begins in a state ψ , and a measurement of the observable f is performed. If the measurement is $\lambda \in \mathbb{R}$, then immediately after the measurement, the system will be in a state ψ' such that

$$\hat{f}\psi'=\lambda\psi'$$
,

where \hat{f} is the self-adjoint operator representation of f.

Remark. Since ψ' is an eigenvector of \hat{f} , this means that a second measurement of f occurring immediately after the first measurement will yield the value λ with probability 1 (see the above proposition about eigenvectors).

The wave function collapse principle only applies to measurements very shortly after the initial measurement; the system will still evolve over time.

Hydrogen Atom Measurements and Uncertainty

The Hamiltonian operator \hat{H} for the hydrogen atom has eigenvalues of the form

$$-\frac{R}{n^2}$$
,

where R is the Rydberg constant and $n \in \mathbb{N}$ denotes a state. The negative value of these eigenvalues is important; they denote that the electron is bound to the nucleus.

If an electron is placed into a state with energy $-\frac{R}{n_1^2}$ (where $n_1 > 1$), it will eventually decay into a state with lower energy, $-\frac{R}{n_2^2}$, with $n_2 < n_1$. The process of decay releases a photon with energy

$$E_{\rm photon} = \frac{R}{n_2^2} - \frac{R}{n_1^2}.$$

The frequency of a photon is proportional to its energy; by observing the frequency of the photon, we can determine the change in energy (and thus the values of n_1 and n_2).

A bound state for the hydrogen atom will be a linear combination of the eigenvectors of \hat{H} with eigenvalues of the aforementioned form.

XIV The spectral decomposition is very cool.

Just as the standard deviation of a random variable Y is found by $\sigma^2 = E(Y^2) - (E(Y))^2$, we can find the *uncertainty* of a self-adjoint operator A, defined as

$$(\Delta_{\psi}A)^2 = \langle A^2 \rangle_{\psi} - (\langle A \rangle_{\psi})^2$$
 ,

recalling that $\langle A \rangle_{\psi} = \langle \psi, A \psi \rangle$.

For a single A, we can find ψ such that $\Delta_{\psi}A < \varepsilon$ for any $\varepsilon > 0$; however, if A and B do not commute, then $\Delta_{\psi}A$ and $\Delta_{\psi}B$ cannot both be made arbitrarily small.

In particular, we know that X and P do not commute; thus yields the famous Heisenberg Uncertainty Principle:

$$(\Delta_{\psi}X)(\Delta_{\psi}P) \geq \frac{\hbar}{2}$$

for all ψ such that $\Delta_{\psi}X$ and $\Delta_{\psi}P$ are defined.

The Schrödinger Equation

In the previous section, we considered the wave function ψ at a fixed time; however, time is not stationary^{XV} and neither ought our system be, meaning we need a way to evolve the quantum system over time.

In a classical system, the Hamiltonian is the generator of time evolution XVI —thus, by the Principle of Representation, there is an operator, \hat{H} , which is the Hamiltonian operator for the system.

We motivated the definition of the momentum operator via the de Broglie hypothesis, which stated that $p = \hbar k$, where k is the angular frequency of the wave function. Similarly, we can look at the relation between energy and the temporal frequency of the wave function,

$$E=\hbar\omega$$
,

to motivate time evolution.

Suppose that ψ_0 has energy E (so ψ_0 is an eigenvector for \hat{H}). Then, the wave function's time dependence should be solely based on the frequency $\omega=E/\hbar$, meaning that if the t=0 state of the system is ψ_0 , the state of the system at any other time t should be

$$\psi(t) = \psi_0 e^{-i\omega t}$$
$$= \psi_0 e^{-it\frac{E}{\hbar}}.$$

This $\psi(t)$ is a function to the differential equation

$$\frac{d\psi}{dt} = -\frac{iE}{\hbar}\psi$$
$$= \frac{E}{i\hbar}\psi$$

with initial value ψ_0 . Finally, rewriting $\hat{H} = E$, we get the Schrödinger equation.

Axiom (Schrödinger Equation). The time evolution of the wave function ψ in a quantum system is generated by

$$\frac{d\psi}{dt} = \frac{1}{i\hbar}\hat{H}\psi,$$

where Ĥ denotes the Hamiltonian operator.

XV Citation needed

XVIThe Hamiltonian is also equal to the total energy of the system.

Proposition (Schrödinger Equation for Operators). Let A be a self-adjoint operator on **H**. Assuming particular domain conditions hold, then

$$\frac{d}{dt}\langle A\rangle_{\psi(t)} = \left\langle \frac{1}{i\hbar} [A, \hat{H}] \right\rangle_{\psi(t)},$$

where $\langle A \rangle_{\psi} = \langle \psi, A \psi \rangle$ and [A, B] = AB - BA.

Remark. Remember that for a function f acting on the classical phase space, $\frac{df}{dt} = \{f, H\}$ acting along a solution to Hamilton's equations.

Proof. Let $\psi(t)$ be a solution to the Schrödinger equation, with $\psi(t) \in \mathsf{Dom}(A) \cap \mathsf{Dom}(\hat{H})$, $A\psi(t) \in \mathsf{Dom}(\hat{H})$, and $H\psi(t) \in \mathsf{Dom}(A)$. Then,

$$\begin{split} \frac{d}{dt} \left\langle \psi(t), A\psi(t) \right\rangle &= \left\langle \frac{d\psi}{dt}, A\psi \right\rangle + \left\langle \psi, A\frac{d\psi}{dt} \right\rangle \\ &= \frac{i}{\hbar} \left\langle \hat{H}\psi, A\psi \right\rangle - \frac{i}{\hbar} \left\langle \psi, A\hat{H}\psi \right\rangle \\ &= \frac{1}{i\hbar} \left\langle \psi, [A, \hat{H}]\psi \right\rangle. \end{split}$$

Remark. If $[A, \hat{H}] = 0$, then nothing interesting happens; thus, for operators to yield time evolution, we need them to not commute with \hat{H} .

For a particle moving in \mathbb{R} , we see that noncommutativity holds for X and P.

If $[A, \hat{H}] = 0$, then we see that $E(A^m)$ is independent of time, effectively meaning that A is conserved.

Proposition (Time Independence of Inner Products of Solutions to the Schrödinger Equation). If $\phi(t)$ and $\psi(t)$ are solutions to the Schrödinger equation, then $\langle \phi(t), \psi(t) \rangle$ and $\|\psi(t)\|$ are independent of time.

Proof. Using the product rule for inner products, we find

$$\begin{split} \frac{d}{dt} \left\langle \phi(t), \psi(t) \right\rangle &= \left\langle \frac{1}{i\hbar} \hat{H} \phi(t), \psi(t) \right\rangle + \left\langle \phi(t), \frac{1}{i\hbar} \hat{H} \psi(t) \right\rangle \\ &= -\frac{1}{i\hbar} \left\langle \hat{H} \phi(t), \psi(t) \right\rangle + \frac{1}{i\hbar} \left\langle \phi(t), \hat{H} \psi(t) \right\rangle \\ &= 0. \end{split}$$

Solving the Schrödinger Equation

The Schrödinger equation is an equation of the form

$$\frac{dv}{dt} = Av$$

for some linear operator A over a Hilbert space. Considering the finite-dimensional Hilbert space \mathbb{C}^n , A is an $n \times n$ matrix. The aforementioned equation thus has the solution

$$v(t) = v_0 e^{At}$$

where v_0 is the initial condition. If A is a diagonalizable matrix, we can calculate e^{At} using the eigenvectors XVII

The Schrödinger equation simply replaces \mathbb{C}^n with \mathbf{H} and A with $\frac{1}{i\hbar}\hat{H}$.

 $^{^{\}mathrm{XVII}}$ In particular, if A is normal, the eigenvectors are the columns of the unitary diagonalization UDU^* , and form an orthonormal basis.

Proposition (Solution to the Schrödinger Equation via Exponentiation). If \hat{H} is a self-adjoint operator on \mathbf{H} , then

$$\psi(t) = \psi_0 e^{-it\hat{H}/\hbar}$$

should be a solution to the Schrödinger equation given that $e^{-it\hat{H}/\hbar}$ is reasonably defined.

Remark. Clearly, XVIII $\psi(t)$ is a solution to the Schrödinger equation.

If \hat{H} is a bounded operator, then $e^{-it\hat{H}/\hbar}$ can be defined by a convergent power series. However, this is rarely the case.

If \hat{H} is unbounded, then we will have to use the spectral theorem — the full spectral theorem will have to wait for a while, but we can examine the case of a point spectrum.

If $\{e_j\}$ is an orthonormal basis for **H** consisting of eigenvectors of \hat{H} with $\hat{H}e_j = \lambda_j e_j$, then we define the exponential by requiring

$$e_i e^{-it\hat{H}/\hbar} = e_i e^{-it\lambda_j/\hbar}$$
.

This construction makes $e^{-it\hat{H}/\hbar}$ a unitary operator, XIX and thus bounded.

It is not necessarily true that every self-adjoint operator A has an orthonormal basis (even self-adjoint ones); nevertheless, the spectral theorem tells us that there is a decomposition of \mathbf{H} into generalized eigenspaces for A.

Unfortunately, this doesn't mean we have *solved* the Schrödinger equation. All it tells us is that $e^{-it\hat{H}/\hbar}$ is bounded, meaning it's defined for all $\psi_0 \in \mathbf{H}$; however, ψ_0 must belong to the domain of \hat{H} (which is not \mathbf{H} , but only a dense subset of it) in order to be a solution.

If ψ is an eigenvector of the Hamiltonian, then the equation

$$\hat{H}\psi = E\psi$$

for some $E \in \mathbb{R}$ is known as the time-independent Schrödinger equation. When we solve the time-independent Schrödinger equation, we are trying to find nonzero values of E and their corresponding ψ that satisfy the equation.

If ψ_0 is a solution to the time-independent Schrödinger equation, then $\psi(t) = \psi_0 e^{-itE/\hbar}$ is the solution to the time-dependent Schrödinger equation with initial condition ψ_0 . Since $\psi(t)$ is a constant multiple of ψ_0 , we say that the solution to the time-independent Schrödinger is a stationary state.

The Schrödinger Equation in $\mathbb R$

Consider a particle moving on the real line. The Hamiltonian for such a particle is written in the form of

$$H(x,p) = \frac{p^2}{2m} + V(x),$$

where V denotes the potential. In the quantum case, we may then consider

$$\hat{H} = \frac{P^2}{2m} + V(X)$$

XVIII My professor loves this phrase.

XIX Notice that $\left|e^{-itE_j/\hbar}\right|=1$.

to be the Hamiltonian operator. Here, V(X) denotes multiplication by V(x), where x is an input to ψ . Thus, we find

$$\hat{H}\psi(x) = -\frac{\hbar^2}{2m}\frac{d^2\psi}{dx^2} + V(x)\psi(x).$$

An operator of the form above is known as the Schrödinger operator; the Hamiltonian operator refers to the operator that generates time-evolution, regardless of its form. The time-dependent Schrödinger equation thus becomes

$$\frac{\partial \psi(x,t)}{\partial t} = \frac{i\hbar}{2m} \frac{\partial^2 \psi(x,t)}{\partial x^2} - \frac{i}{\hbar} V(x) \psi(x,t).$$

The time-dependent Schrödinger equation is thus a linear partial differential equation. XX For a particle on \mathbb{R} , the time-independent Schrödinger equation is a linear ordinary differential equation with nonconstant coefficients. XXI

Expected Position and Momentum for Solutions to the Schrödinger Equation

Proposition (Time Evolution of Position and Momentum). Let $\psi(t)$ be a solution to the time-dependent Schrödinger equation with sufficiently nice V and sufficiently nice $\psi_0 = \psi(0)$. Then,

$$\begin{split} \frac{d}{dt} \langle X \rangle_{\psi(t)} &= \frac{1}{m} \langle P \rangle_{\psi(t)} \\ \frac{d}{dt} \langle P \rangle_{\psi(t)} &= - \langle V'(X) \rangle_{\psi(t)} \,. \end{split}$$

Remark. We require these assumptions to ensure that \hat{H} is a self-adjoint operator and that the domain conditions hold.

Note that $-\langle V'(X)\rangle_{\psi(t)}\neq -V'\left(\langle X\rangle_{\psi(t)}\right)$. This solution to the Schrödinger equation is not the quantum analogue to the solutions of Hamilton's equations.

Although expected position and expected momentum do not exactly follow the classical trajectories, if $\psi(x)$ is very closely concentrated about $x = x_0$, then the particle will approximately follow the classical trajectory.

The Heisenberg Approach

The previous approach (known as the Schrödinger approach) towards analyzing the quantum system comes from understanding how the wave functions evolve over time, with the operators remaining stationary. We will now examine an approach where the operators change over time with the states remaining stationary.

In the Heisenberg approach, every self-adjoint A evolves in time according to the operator differential equation X^{XXII}

$$\frac{dA(t)}{dt} = \frac{1}{i\hbar}[A(t), \hat{H}].$$

Note that since \hat{H} commutes with itself, it remains constant in time — this is akin to the classical Hamiltonian remaining constant along a solution to Hamilton's equations.

XXUnfortunately I haven't taken PDEs yet, so I can't really explain how to solve this.

XXI Neither have I taken ODEs, but I will be taking it in Fall 2024!

XXIII will call this the "modified Schrödinger equation."

The spectral theorem provides us a way to construct a family of unitary operators $e^{-it\hat{H}/\hbar}$ that computes time-evolution of states in the Schrödinger approach; similarly, we can find that the solutions in the operator differential equation are of the form

$$A(t) = Ae^{-it\hat{H}/\hbar}e^{it\hat{H}/\hbar}.$$

If ψ is the state of the system, then the expectation of A(t) in the state is defined to be $\langle \psi, A(t)\psi \rangle$; then

$$\langle \psi, A(t)\psi \rangle = \left\langle \psi, e^{it\hat{H}/\hbar} A e^{-it\hat{H}/\hbar} \psi \right\rangle$$
$$= \left\langle e^{-it\hat{H}/\hbar} \psi, A e^{-it\hat{H}/\hbar} \psi \right\rangle$$
$$= \left\langle \psi(t), A \psi(t) \right\rangle,$$

which is the time-evolved state of the system. Since $\langle \psi(t), A\psi(t) \rangle$ is the expectation of the value of A in the state $\psi(t)$, it must be the case that the Heisenberg approach and the Schrödinger approach are equivalent in their physics.

Proposition (Hamiltonian in the Heisenberg Approach). Let $\hat{H} = \frac{P^2}{2m} + V(X)$, where V is a bounded below polynomial. Then, for any $t \in \mathbb{R}$,

$$\hat{H} = \frac{1}{2m} (P(t))^2 + V(X(t)).$$

Remark. Notice that the Hamiltonian is independent of time (since $[\hat{H}, \hat{H}] = 0$), even though P and X depend on time.

Lemma (Moments of Solutions). Suppose A is a self-adjoint operator on \mathbf{H} and A(t) is a solution to the modified Schrödinger equation with A(0) = A. Then, $(A(t))^m$ is also a solution to the modified Schrödinger equation. In other words, the time evolution of the mth power of A is the same as the mth power of the time evolution of A.

Proof.

$$e^{it\hat{H}/\hbar}A^{m}e^{-it\hat{H}/\hbar} = \prod_{k=1}^{m} \left(e^{it\hat{H}/\hbar}Ae^{-it\hat{H}/\hbar}\right)$$
$$= \left(e^{it\hat{H}/\hbar}Ae^{-it\hat{H}/\hbar}\right).$$

Proposition (Analogue to Hamilton's Equations). Suppose $\hat{H} = \frac{(P(t))^2}{2m} + V(X(t))$. Then,

$$\frac{dX}{dt} = \frac{1}{m}P(t)$$
$$\frac{dP}{dt} = -V'(X(t)).$$

Remark. Notice that it is the functions X(t) and P(t) that satisfy the analogue to Hamilton's equations; the expectation values satisfy

$$\frac{d}{dt} \langle X(t) \rangle_{\psi} = \frac{1}{m} \langle P(t) \rangle_{\psi}$$
$$\frac{d}{dt} \langle P(t) \rangle_{\psi} = -\langle V'(X(t)) \rangle_{\psi}.$$

Particle in a Box

We want to solve the time-independent Schrödinger equation for the case of a particle that is confined to move in the interval $x \in [0, L]$. That is, we want to find all the eigenvectors and eigenvalues of $\hat{H}\psi = E\psi$.

The particle has potential 0 for $x \in [0, L]$ and very large C outside [0, L].