1. States, Observables, and Symmetries

The basic ingredients of a theory of a physical system are mathematical objects describing 1) a set $S$ of states of the system, 2) the set $O$ of observables of a system, that is, real-valued properties of the system which can be experimentally determined, and 3) the group $G$ of symmetries of the system. The difference between classical physics and quantum theory arises from very different choices of mathematical objects used to describe states, observables, and symmetries. But in both classical and quantum physics certain mathematical structures appear that are associated with $(S,O,G)$. Three of the most important are:

1. Given a state $\psi \in S$ and an observable $A \in O$ there is a probability measure $\mu_{A,\psi}$ on $\mathbb{R}$ which represents the probability distribution of values that the observable $A$ may be measured to have when the system is in the state $\psi$. (We are not assuming that the observable $A$ has a unique numerical value in a given state $\psi$; repeated measurements of the same observable a system in the same state may have different values, and these values are distributed according to a probability distribution.)

A probability measure $\mu$ on a measure space $X$ is a finite measure such that $\mu(X) = 1$. Measures on a topological space, e.g. $\mathbb{R}$, will always be assumed to be defined on the $\sigma$-algebra of Borel sets.

2. The group $G$ has an action $U$ on the set $S$ of states.

An action $U$ of a group $G$ on a set $S$ is a mapping $U: G \times S \to S$, usually written as $(g,x) \mapsto U(g)x$, such that: 1) $U(e)x = x$ where $e$ is the identity in $G$; 2) $U(gh)x = U(g)(U(h)x)$. These imply $U(g^{-1}) = U(g)^{-1}$.

The most important and simple example is where $G = \mathbb{R}$ and $U$ is time evolution; i.e. given a state $\psi \in S$ and $t \in \mathbb{R}$, the state $U(t)\psi$ is obtained by simply waiting $t$ units of time, letting the state evolve according to the physical laws the system obeys. The requirement $U(t+s) = U(t)U(s)$ amounts to saying that the physical laws don't change with time, i.e. there are no time-dependent external influences on the system.

3. The group $G$ is a topological group, that is, a group with a topology such multiplication is continuous from $G \times G$ to $G$ and the map $g \mapsto g^{-1}$ is continuous from $G$ to $G$. A continuous homomorphism $t \mapsto g_t$ from $\mathbb{R}$ to $G$ is called a one-parameter subgroup,
though strictly speaking its range is a subgroup of $G$. Given a one-parameter subgroup let $U(t) = U(g_t)$. Then associated to the one-parameter subgroup there is an observable $A \in O$ called the associated conserved quantity, which has the property that $\mu_{A,U(t)} \psi$ is independent of $t$. For example, in physics the conserved quantity associated with time evolution is called energy or the Hamiltonian. In systems which have symmetry under spatial translations, the associated conserved quantities are called momenta (in the $x, y$, and $z$ directions, for example). In systems which have rotational symmetry the associated conserved quantities are called angular momenta (about the $x, y$ and $z$ axes, for example). In electromagnetism there is a one-parameter group of symmetries called global gauge transformations which corresponds to the conserved quantity electric charge.

This relation between states, observables, and symmetries is less obvious and more interesting than the first two; it often goes by the name of Nöther’s Theorem, after Emmy Nöther, the student of Hilbert, since a version of it applicable to problems in variational calculus was proved in her ‘Invariante Variationsprobleme’ article in 1918, available in English in Transport Theory and Stat. Phys. 1, 186-207 (1971).

We now will say what mathematical objects fill the roles of $(S, O, G)$ in classical mechanics and quantum theory. In the case of quantum theory we will immediately be led to the spectral theorem and Stone’s theorem, which we will then prove.

2. States, Observables, and Symmetries in Classical Mechanics

In classical mechanics, the set of states $S$ is $\mathbb{R}^{2n} = \mathbb{R}^n \oplus \mathbb{R}^n$, called phase space, and the first summand $\mathbb{R}^n$ is called configuration space. One uses coordinates $\{q_i\}_{i=1}^n$ for configuration space; these coordinates are called positions, and one uses coordinates $\{q_i, p_i\}_{i=1}^n$ for phase space; the $p_i$ being called momenta.

The relation between these ‘momenta’ and the ‘momenta’ mentioned above is indirect. A typical example of phase space appears in the description of a system of $n$ point particles in ordinary 3-dimensional space; each particle’s state is described by its position $\vec{q}_i = (q_{1i}, q_{2i}, q_{3i})$ and its momentum, or mass times velocity vector, $\vec{p}_i = (p_{1i}, p_{2i}, p_{3i})$, and the whole assemblage is described by the following vector in the phase space $S = \mathbb{R}^{3n} \oplus \mathbb{R}^{3n}$: $(\vec{q}_1, \ldots, \vec{q}_n, \vec{p}_1, \ldots, \vec{p}_n)$.

The set of observables $O$ is the space $C(\mathbb{R}^{2n})$ of continuous real-valued functions on phase space. Given a state $\psi \in \mathbb{R}^{2n}$ and an observable $A \in C(\mathbb{R}^{2n})$, the probability measure $\mu_{A,\psi}$ mentioned above is just $\delta_{A(\psi)}$, the Dirac delta measure at $A(\psi) \in \mathbb{R}$. In other words, if the state is $\psi$ the value of any observable $A$ is always measured to be $A(\psi)$; there’s no ‘uncertainty’ in the value of $A$.

We will consider only the simplest case, in which the group $G$ of symmetries is just $\mathbb{R}$. Given a smooth function $H \in C^\infty(\mathbb{R}^{2n})$, we define the associated vector field $v_H$ on
A vector field \( v \) on \( \mathbb{R}^{2n} \) is **integrable** if there is an action \( U \) of \( \mathbb{R} \) on \( \mathbb{R}^{2n} \) such that

\[
\frac{d}{dt} U(t) \psi = v(U(t) \psi) .
\]

If \( v_H \) is integrable and \( U \) is an action of \( \mathbb{R} \) on \( \mathbb{R}^{2n} \) such that \( \frac{d}{dt} U(t) \psi = v_H(U(t) \psi) \), and we write in coordinates \( U(t) \psi = (q_i, p_i) \in \mathbb{R}^{2n} \), then **Hamilton’s equations** hold:

\[
\frac{dq_i}{dt} = \frac{\partial H}{\partial p_i} ; \quad \frac{dp_i}{dt} = -\frac{\partial H}{\partial q_i} .
\]

In classical mechanics the action \( U \) of the group \( \mathbb{R} \) on states is required to be one arising from some \( H \in C^\infty(\mathbb{R}^{2n}) \) in this way, and \( H \) is the conserved quantity associated to \( U \).

**Exercise:** Show that Hamilton’s equations are indeed equivalent to

\[
\frac{d}{dt} U(t) \psi = v_H(U(t) \psi) .
\]

Show that if \( H \in C^\infty(\mathbb{R}^{2n}) \) and \( v_H \) is integrable, the directional derivative \( v_H H \) equals zero. Show that this implies that the measure \( \mu_{H,U(t)} \) is independent of \( t \) for any state \( x \in \mathbb{R}^{2n} \). (This is a baby version of Noether’s theorem.)

**Exercise:** Show that if \( H = \frac{1}{2} \sum_{i=1}^{n} q_i^2 + p_i^2 \) (the **harmonic oscillator Hamiltonian**) the vector field \( v_H \) is integrable by explicitly solving Hamilton’s equations. Show that if \( H = \sum_{i=1}^{n} \frac{1}{2}(q_i^2 + p_i^2) + \frac{1}{2} \lambda_i p_i^4 \) (the **anharmonic oscillator Hamiltonian**) the vector field \( v_H \) is integrable if and only if all the **coupling constants** \( \lambda_i \) are non-negative.

The setup for classical mechanics as we have described it can be substantially generalized. The three most important directions of generalization are: 1) The group of symmetries \( G \) can be an arbitrary **Lie group** (a group which is also a smooth manifold, such that the group operations are smooth maps). 2) The phase space \( S \) can be an arbitrary **symplectic manifold** (a manifold with a non-degenerate closed 2-form). 3) The group \( G \) and the phase space \( S \) can be infinite-dimensional; this occurs in the study of classical field theory, i.e. the study of wave equations.
3. States, Observables and Symmetries in Quantum Theory

In quantum mechanics the set of states \( S \) is the unit sphere in a Hilbert space \( H \); i.e.

\[
S = \{ \psi \in H : \| \psi \| = 1 \}.
\]

As we shall see, for many purposes it’s better to take the state space to be the quotient of this sphere by the equivalence relation: \( \psi \sim \phi \) if \( \psi = c\phi \) for some \( c \in \mathbb{C} \), i.e. our present description is ‘redundant’, but this is unimportant for introductory purposes.

The set of observables \( O \) is then the set of self-adjoint operators on \( H \). Recall that a operator on \( H \) is a linear map \( A \) from a linear subspace of \( H \) to \( H \). This subspace, the domain of \( A \), is denoted \( D(A) \), and we are not requiring even that it be dense in \( H \). If \( D(A) \) is dense in \( H \) we say that \( A \) is densely defined. An operator with domain equal to all of \( H \), satisfying

\[
\exists c \geq 0 \forall \psi \in H \quad \| A\psi \| \leq c\| \psi \|
\]

is called a bounded operator, and the space of all bounded operators on \( H \) is denoted by \( \mathcal{L}(H) \). The adjoint of an operator \( A \) is supposed to have the property \( \langle \phi, A\psi \rangle = \langle A^*\phi, \psi \rangle \), where we use \( \langle \cdot, \cdot \rangle \) for the inner product in \( H \). (Also, we follow the physicists’ convention that the inner product is linear in the second variable.) But for operators that are not bounded this requires some care in formulation. Thus we define the adjoint of a densely defined operator \( A \) to be the operator \( A^* \) with domain

\[
D(A^*) = \{ \phi \in H : \exists \phi' \in H \forall \psi \in D(A) \quad \langle \phi, A\psi \rangle = \langle \phi', \psi \rangle \}
\]

and noting that if \( \phi' \) with the above property exists it is unique (because \( A \) is densely defined), we define \( A^*\phi \) for \( \phi \in D(A^*) \) to be the unique \( \phi' \) such that

\[
\forall \psi \in D(A) \quad \langle \phi, A\psi \rangle = \langle \phi', \psi \rangle.
\]

An operator \( A \) on a Hilbert space is said to be self-adjoint if it is densely defined and equal to its adjoint (which in particular means that \( D(A) = D(A^*) \)).

The definition of adjoint may look complicated, but it is the most natural possible definition for densely defined operators. It foreshadows the fact that densely defined (unbounded) operators are altogether more subtle than bounded operators. Physicists typically ignore these nuances, and indeed the fastest way to tell if a physicist is a ‘mathematical physicist’ is to ask him to define ‘adjoint’.

Given a state \( \psi \) in \( S \) and a self-adjoint operator \( A \) on \( H \) there is a probability measure \( \mu_{A,\psi} \) on \( \mathbb{R} \), called the ‘spectral measure’, which describes the probability distribution
of measured values for $A$ when the system is in the state $\psi$. This measure arises as a
consequence of the spectral theorem, which we shall soon state and prove.

A bounded operator $T$ on a Hilbert space $H$ is \textbf{unitary} if $TT^* = T^*T = I$, or
equivalently if $T$ is onto and $\|T\psi\| = \|\psi\|$ for all $\psi \in H$. Recall that an action $U$ of a
group $G$ on a vector space $V$ is a \textbf{representation} if $U(g)$ is linear for all $g \in G$; if in
addition $S$ is a Hilbert space and $U(g)$ is unitary for all $g \in G$ we say $U$ is a \textbf{unitary
representation}. If $G$ is a topological group we say that a representation of $G$ on a
Hilbert space is \textbf{strongly continuous} if $g_\alpha \to g$ in $G$ implies $\|U(g_\alpha)\psi - U(g)\psi\| \to 0$
for all $\psi \in H$. In quantum theory one requires that the group $G$ of symmetries be a
topological group and that the action $U$ of $G$ on the states $S$ come from a strongly continuous
unitary representation of $\mathbb{R}$ on $H$. Stone’s theorem then associates to any one-parameter subgroup
$g_t$ of $G$ a conserved quantity $A$, a self-adjoint operator on $H$ also called the ‘generator’ of
$U(t) = U(g_t)$.

We recall here the three most important topologies on $\mathcal{L}(H)$. The \textbf{norm} topology is defined
by: $T_\alpha \to T$ if $\|T_\alpha - T\| \to 0$. The \textbf{strong} topology is defined by: $T_\alpha \to T$ if $\|T_\alpha \psi - T\psi\| \to 0$
for all $\psi \in H$. And the \textbf{weak} topology is defined by: $T_\alpha \to T$ if $\langle \psi, (T_\alpha - T)\psi \rangle \to 0$
for all $\psi, \phi \in H$. A function from or to $\mathcal{L}(H)$ is said to be norm-continuous, strongly continuous, or
weakly continuous depending on which topology is used for $\mathcal{L}(H)$.

Exercise - Let $\ell^2$ denote the Hilbert space of complex sequences $\psi = (\psi_1, \psi_2, \ldots)$ with
finite norm $\|\psi\| = (\sum |\psi_i|^2)^{1/2}$, given the inner product $\langle \psi, \phi \rangle = \sum \overline{\psi}_i \phi_i$. Let $T$ be the
\textbf{left shift} operator, given by $T(\psi_1, \psi_2, \ldots) = (\psi_2, \psi_3, \ldots)$. Calculate $T^*(\psi_1, \psi_2, \ldots)$ and
explain why $T^*$ is called the \textbf{right shift} operator. These shift operators illustrate some
of the things that bounded linear operators can do in infinite dimensions that they cannot
in finite dimensions. For example, show that $T^*$ is an \textbf{isometry}, e.g. $\|T^*\psi\| = \|\psi\|$, but
not unitary. Show that an isometry of a finite-dimensional Hilbert space is unitary. Show
that $T$ and $T^*$ do not commute.

Exercise - Show that norm convergence implies strong convergence, which implies
weak convergence, but that the converses do not hold. (For counterexamples, consider the
sequences $\{T^n\}$ and $\{T^{*n}\}$, where $T$ is the left shift.) Show that converses hold in a
finite-dimensional Hilbert space. Show also that if $U_\alpha$ are \textbf{unitary} and $U_\alpha \to U$ weakly
then $U_\alpha \to U$ strongly; thus a weakly continuous unitary representation is automatically
strongly continuous.

Exercise - Show that for bounded operators $S,T$ we have $T^{**} = T$, $\|T^*\| = \|T\|$, $\|T^*T\| = \|T\|^2$, $(cT)^* = \overline{c}T^*$, $(S+T)^* = S^* + T^*$, and $(ST)^* = T^* S^*$. Show that the map
$T \mapsto T^*$ is norm-continuous and weakly continuous but not strongly continuous. (For the
counterexample consider shift operators.)
4. The Spectral Theorem - Bounded Case

We begin by proving the spectral theorem for bounded self-adjoint operators and then deal with the unbounded case. First we define the spectrum of a bounded operator (not necessarily self-adjoint) $T \in \mathcal{L}(H)$ to be the set of $\lambda \in \mathbb{C}$ such that $T - \lambda I$ is not invertible, i.e. is not a bijection with a bounded inverse. The spectrum of $T$ is denoted $\sigma(T)$. The complement of the spectrum of $T$ is called the resolvent set of $T$ and denoted $\rho(T)$. The exercises below show that the spectrum is a natural generalization of the set of eigenvalues of a matrix.

Exercise - Show that if $H = \mathbb{C}^n$ the spectrum of $T \in \mathcal{L}(H)$ is just the set of $\lambda \in \mathbb{C}$ such that $\det(T - \lambda I) = 0$. Show in this case that the spectrum of $T$ is the set of eigenvalues of $T$ (for self-adjoint $T$ this is easy; for general $T$ one can use the Jordan canonical form - see any advanced book on linear algebra, e.g. Birkhoff and MacLane, *A Survey of Modern Algebra*.)

Exercise - Given a bounded operator $T$ on $H$ define $\lambda \in \mathbb{C}$ to be an approximate eigenvalue of $T$ if for any $\epsilon > 0$ there exists $\psi \in H$ with $\|\psi\| = 1$ and $\|T\psi - \lambda \psi\| \leq \epsilon$. Show that if $\lambda$ is an approximate eigenvalue for $T$ then $\lambda \in \sigma(T)$. For a partial converse, show that if $A$ is a bounded self-adjoint operator on a Hilbert space and $\lambda \in \sigma(A)$, then $\lambda$ is an approximate eigenvalue of $A$. Hint for the second part: show that for any densely defined operator $T$ on $L(H)$,

$$\text{Ran}(T) = \text{Ker}(T^*)$$

so that for bounded self-adjoint $A$, $(A - \lambda I)$ is invertible unless $\lambda$ is an approximate eigenvalue of $A$.

Exercise - Show that $\lambda \in \sigma(T)$ implies $\overline{\lambda} \in \sigma(T^*)$. Show that if $A$ is bounded and self-adjoint $\sigma(A) \subseteq \mathbb{R}$.

The Spectral Mapping Theorem (for polynomials). *If $T$ is a bounded operator and $P$ is a nonconstant polynomial then $\sigma(P(T)) = \{P(\lambda) : \lambda \in \sigma(T)\}$. (For short we simply write $\sigma(P(T)) = P(\sigma(T))$.)*

Proof - Let $\lambda \in \sigma(T)$. Since $x = \lambda$ is a root of $P(x) - P(\lambda)$, we can factor:

$$P(x) - P(\lambda) = (x - \lambda)Q(x),$$

where $Q(x)$ is a polynomial with $\lambda$ as a root. Since $T$ is self-adjoint, $T - \lambda I$ is invertible on $\text{Ran}(T)$, and hence invertible on $\mathbb{C}^n$. Therefore, $P(T)$ is invertible on $\mathbb{C}^n$, and $\lambda \in \sigma(P(T))$. Conversely, if $\lambda \in \sigma(P(T))$, then $P(T) - P(\lambda)$ is not invertible, and hence $T - \lambda I$ is not invertible on $\mathbb{C}^n$. Therefore, $\lambda \in \sigma(T)$.
so $P(T) - P(\lambda) = (T - \lambda I)Q(T)$. Since the operator $T - \lambda I$ is not invertible neither is $P(T) - P(\lambda)$, i.e. $P(\lambda) \in \sigma(P(T))$.

Conversely suppose $\mu \in \sigma(P(T))$. Factor $P(x) - \mu$ as $a(x - \lambda_1) \cdots (x - \lambda_n)$. Then $P(T) - \mu I = a(T - \lambda_1) \cdots (T - \lambda_n)$. If all the $\lambda_i$ are not in $\sigma(T)$ then we have

$$(P(T) - \mu I)^{-1} = a^{-1}(T - \lambda_1)^{-1} \cdots (T - \lambda_n)^{-1},$$

so $\mu$ is not in $\sigma(P(T))$, a contradiction. Thus there must be some $\lambda_i \in \sigma(T)$, and

$$P(\lambda_i) - \mu = a(\lambda_i - \lambda_1) \cdots (\lambda_i - \lambda_n) = 0$$

so $\mu = P(\lambda_i)$. $\square$

**Theorem.** The spectral radius of a bounded operator $T$ is defined to be $r(T) = \sup\{|\lambda|: \lambda \in \sigma(T)\}$. If $A$ is a bounded self-adjoint operator on $H$ then $r(A) = \|A\|$.

**Proof** - If $\lambda \in \sigma(A)$ then $\lambda$ is an approximate eigenvalue for $A$ so there exists $\psi_n \in H$ with $\|\psi_n\| = 1$ and $\|A\psi_n\| \to |\lambda|$. Thus $\|A\| \geq |\lambda|$, so $\|A\| \geq r(A)$.

Now choose $\psi_n$ such that $\|\psi_n\| = 1$ and $\|A\psi_n\| \to \|A\|$. Let $\lambda = \|A\|$. If $A$ is self-adjoint then

$$\|A^2\psi_n - \lambda^2\psi_n\|^2 = \langle A^2\psi_n - \lambda^2\psi_n, A^2\psi_n - \lambda^2\psi_n \rangle$$

$$= \|A^2\psi_n\|^2 - 2\lambda^2\|A\psi_n\|^2 + \lambda^4$$

$$\leq (\|A\| \|A\psi_n\|)^2 - 2\lambda^2\|A\psi_n\|^2 + \lambda^4$$

$$= -\lambda^2\|A\psi_n\|^2 + \lambda^4 \to 0.$$

Thus $\lambda^2$ is an approximate eigenvalue for $A^2$, so $\lambda^2 \in \sigma(A^2) = \sigma(A)^2$. It follows that $\|A\| \leq r(A)$. Thus $\|A\| = r(A)$. $\square$

The remarkable thing about this theorem is that it relates a topological concept, the norm, to a concept which is defined purely algebraically, the spectral radius. This interplay of topology and algebra is behind the beautiful theory of $C^*$-algebras, a kind of algebra over the complex numbers which has a norm and $*$ operator satisfying axioms like those which hold for $\mathcal{L}(H)$. $C^*$-algebras are important in the mathematics of quantum field theory, and in understanding the relation between classical and quantum physics.
Exercise - Fill out the following proof of a generalization of the theorem above: if \( T \) is a bounded operator then \( r(T) \leq \|T\| \) and
\[
r(T) = \limsup_{n \to \infty} \|T^n\|^{1/n}.
\]
The idea is to use the Laurent series for \((T - \lambda I)^{-1}\), which is modeled after the series for \((z - \lambda)^{-1}\):
\[
(T - \lambda I)^{-1} = -\sum_{n=0}^{\infty} \lambda^{-n-1} T^n.
\]
In the region in \( \mathbb{C} \) where the series is norm-convergent, term-by-term multiplication is justified in showing that \((T - \lambda I)^{-1}\) multiplied by \(-\sum_{n=0}^{\infty} \lambda^{-n-1} T^n\), in either order, is the identity. Recall however Hadamard’s theorem that the radius of convergence of \( \sum a_n z^n \) is \( R = \frac{1}{\limsup_{n \to \infty} |a_n|^{1/n}} \), and that \( |\sum a_n z^n| \to \infty \) for some sequence of \( z \)'s with \( |z| \to R \). The same argument (fortified with the uniform boundedness theorem) shows that the radius of convergence of \( \sum T^n z^n \) is \( r = \frac{1}{\lim\sup\|T^n\|^{1/n}} \), and that \( \|\sum T^n z^n\| \to \infty \) for some sequence of \( z \)'s with \( |z| \to R \). Taking \( \lambda = 1/z \), it follows that \((T - \lambda)^{-1}\) exists for \( \lambda \geq R = \lim\sup\|T^n\|^{1/n} \), and
\[
\|(T - \lambda I)^{-1}\| \to \infty
\]
for some sequence of \( \lambda \)'s with \( |\lambda| \) converging to \( R \). This immediately implies that \( r(T) \leq R \), and to show that it also implies \( r(T) = R \) it suffices to show that \( \|(T - \lambda I)^{-1}\| \) is continuous on \( \rho(T) \), since then \( r(T) < R \) would imply have \( \|(T - \lambda I)^{-1}\| \) bounded near \( |\lambda| = R \), a contradiction. The continuity follows from a similar power series argument.

Exercise - Let \( T \) denote the left shift operator on \( \ell^2 \). Compute the spectral radius of \( T \), show that the whole unit circle in \( \mathbb{C} \) consists of approximate eigenvalues of \( T \), and show that \( 0 \) is an eigenvalue of \( T \). Show that the whole unit circle in \( \mathbb{C} \) consists of approximate eigenvalues of \( T^* \) and that \( 0 \) is in the spectrum of \( T^* \), but \( 0 \) is not an approximate eigenvalue of \( T^* \).

**Theorem.** If \( T \) is a bounded operator its spectrum is compact.

Proof - In the previous theorem we showed \( r(T) \leq \|T\| \), so \( \sigma(T) \) is bounded and it suffices to show that \( \sigma(T) \) is closed or equivalently that the resolvent set \( \rho(T) \) is open. Suppose \( \lambda_0 \in \rho(T) \) and let \( R = (T - \lambda I)^{-1} \). Then for \( \lambda \) sufficiently close to \( \lambda_0 \) the power series
\[
\sum_{n=0}^{\infty} (\lambda - \lambda_0)^n R^{n+1}
\]
is norm-convergent and
\[(T - \lambda I)(\sum_{n=0}^{\infty} (\lambda - \lambda_0)^n R^{n+1}) = \sum_{n=0}^{\infty} [(T - \lambda_0 I)(\lambda - \lambda_0)^n R^{n+1} + (\lambda_0 - \lambda)(\lambda - \lambda_0)^n R^{n+1}] = \sum_{n=0}^{\infty}[(\lambda - \lambda_0)^n R^n - (\lambda - \lambda_0)^n R^{n+1}] = I .\]
The same argument shows that
\[\left(\sum_{n=0}^{\infty} (\lambda - \lambda_0)^n R^{n+1}\right)(T - \lambda I) = I,\]
so we have found an explicit inverse for \((T - \lambda I)\). Thus \(\lambda \in \rho(T)\). \(\square\)

Note the use of power series techniques in both the proof of this theorem and in the proof that \(r(T) = \lim \sup_{n \to \infty} \|T^n\|^{1/n}\) in the exercise above. Such techniques from complex analysis are very useful in operator theory; \((T - \lambda I)^{-1}\) is called an analytic operator-valued function on \(\rho(T)\) because it admits such local power series expansions in \(\lambda\).

**Spectral Theorem - continuous functional calculus.** If \(A\) is bounded self-adjoint operator on the Hilbert space \(H\) there is a unique map \(\pi: C(\sigma(A)) \to \mathcal{L}(H)\) such that:

1. \(\pi\) is a \(\ast\)-homomorphism, i.e.
   \[\pi(\lambda f) = \lambda \pi(f) , \quad \pi(f + g) = \pi(f) + \pi(g) , \quad \pi(fg) = \pi(f) \pi(g) , \quad \pi(f^*) = \pi(f)^* , \quad \pi(1) = I ,\]
2. \(\pi\) is continuous, i.e., \(\|\pi(f)\| \leq C\|f\|_\infty\).
3. If \(h\) is the function \(h(\lambda) = \lambda\), then \(\pi(h) = A\).

We shall usually write \(f(A)\) for \(\pi(f)\). In this notation, \(\pi\) has the following properties:

4. If \(f \geq 0\) then \(f(A) \geq 0\).
5. \(\|f(A)\| = \|f\|_\infty\).
6. \(\sigma(f(A)) = f(\sigma(A))\). (The spectral mapping theorem.)
7. If $A\psi = \lambda \psi$ then $f(A)\psi = f(\lambda)\psi$.

Proof - We first define $\pi$ on polynomials by $\pi(p) = p(A)$. The idea is to show that the definition of a polynomial function of a self-adjoint operator extends to continuous functions, obtaining a ‘continuous functional calculus’. The trick is to use Weierstrass’ theorem, which says that the set of polynomials is dense in $C(\sigma(A))$, since $\sigma(A)$ is compact. If we can show 2 for polynomial $f$, it will then follow that $\pi$ extends uniquely from the polynomials to all of $C(\sigma(A))$, and satisfies 2 for all $f$. (Here we are using the easy but all-important B.L.T. theorem (for ‘bounded linear transformation’), which says that if $X$ and $Y$ are Banach spaces, $D$ is a dense subspace of $X$, and $f:D \to Y$ satisfies $\|f(x)\| \leq C\|x\|$ for all $x \in D$, then $f$ extends uniquely to a bounded linear function from $X$ to $Y$.) Thus we first verify 2 for a polynomial $p$. We have

$$
\|\pi(f)\| = \|p(A)\|
= (\|p(A)^*p(A)\|)^{1/2}
= \|p(A)\|^2
= r((pA)(A))^{1/2},
$$

the latter since $p(A)^*p(A) = (pA)(A)$ is self-adjoint. By the spectral mapping theorem for polynomials we have

$$
r((pA)(A)) = \sup\{p(\lambda)\lambda: \lambda \in \sigma(A)\}
= \|p\|_\infty^2,
$$

Thus

$$
\|\pi(p)\| = \|p\|_\infty
$$

if $p$ is the restriction of a polynomial to $\sigma(A)$, verifying 2 for polynomials.

It follows from the B.L.T. theorem that we may uniquely extend $\pi$ to continuous functions, and that 2 and 7 hold. Clearly 3 holds, and since 1 holds for polynomial $f,g$ its easy to see that it extends by continuity to hold for all $f,g \in C(\sigma(A))$. It also follows from the B.L.T. theorem that $\pi$ as we’ve defined it is the unique function satisfying 1 - 3.

Property 4 clearly holds for polynomial $f$ and extends by continuity to all $f \in C(\sigma(A))$. We’ve shown that 5 holds. 7 holds for polynomials and extends by continuity to all $f \in C(\sigma(A))$.

It remains to prove 6, which we know holds for polynomials. We show that $f(\sigma(A)) \subseteq \sigma(f(A))$ as follows. Let $\lambda \in \sigma(A)$; $\lambda$ is an approximate eigenvalue of $A$ so choose $\psi_n \in H$ with $\|\psi_n\| = 1$ and $\|(A - \lambda I)\psi_n\| \leq 1/n$. Choose polynomials $P_m$ converging uniformly to
\( f \) on \( \sigma(A) \), with \( \|P_m - f\| \leq 1/m \). Then for any \( m \),
\[
\|(f(A) - f(\lambda)I)\psi_n\| \leq \|(f(A) - P_m(A))\psi_n\| + \|((P_m(A) - P_m(\lambda)I)\psi_n\|
\leq \|f(A) - P_m(A)\| + \|((P_m(A) - P_m(\lambda))\psi_n\| + \|P_m - f\|_{L^\infty(\sigma(A))}
\leq \|(P_m(A) - P_m(\lambda))\psi_n\| + 2/m.
\]
Since \( \|\psi_n\| = 1 \) and \( \|A\psi - \lambda\psi\| \leq 1/n \), we have, for any \( k \geq 0 \),
\[
\|A^k\psi_n - \lambda^k\psi_n\| = \|(A^{k-1} + \lambda A^{k-2} + \cdots + \lambda^{k-1})(A - \lambda I)\psi_n\|
\leq \|A^{k-1} + \lambda A^{k-2} + \cdots + \lambda^{k-1}\|/n,
\]
which goes to zero as \( n \to \infty \), so
\[
\lim_{n \to \infty} \|(P_m(A) - P_m(\lambda))\psi_n\| = 0,
\]
so for \( n \) large enough this implies that
\[
\|(f(A) - f(\lambda)I)\psi_n\| \leq 3/m.
\]
Since \( m \) is arbitrary this implies that
\[
\|(f(A) - f(\lambda)I)\psi_n\| \to 0.
\]
Thus \( f(\lambda) \) is an approximate eigenvalue for \( f(A) \), so \( f(\lambda) \in \sigma(f(A)) \).

To show that \( \sigma(f(A)) \subseteq f(\sigma(A)) \), suppose that \( \lambda \notin f(\sigma(A)) \). Then the function
\[
g(x) = (f(x) - \lambda)^{-1}
\]
is continuous on \( \sigma(A) \), so we obtain \( g(A) \) by the functional calculus, and
\[
g(A)(f(A) - \lambda I) = I = (f(A) - \lambda I)g(A)
\]
by the \( \ast \)-homomorphism property. Thus \( f(A) - \lambda I \) is invertible, so \( \lambda \notin \sigma(f(A)) \). \( \square \)

The spectral theorem has many versions, and a long and tangled history. There are three
main forms of the spectral theorem: the functional calculus form, the projection-valued measure
form (see the remarks on spectral projections in the next section), and the following form: if
\( A \) is a self-adjoint operator on a Hilbert space \( H \), there is a measure space \( (X, \mu) \), a real-
valued measurable function \( f \) on \( X \), and a unitary transformation \( U : H \to L^2(X) \) such that
\( A = U^{-1}M_fU \), where \( M_f \) denotes multiplication by \( f \). That is, every self-adjoint operator is
unitarily equivalent to a multiplication operator.

Exercise - Prove the B.L.T. theorem.
Note that if $\psi \in H$ the above theorem implies that $f \mapsto \langle \psi, f(A)\psi \rangle$ is a continuous linear functional from $C(\sigma(A))$ to $\mathbb{C}$, which moreover is positive by property 4. Thus by the Riesz-Markov theorem there is a unique measure $\mu_{A,\psi}$ on $\sigma(A)$, the spectral measure, such that

$$\langle \psi, f(A)\psi \rangle = \int_{\sigma(A)} f(\lambda) d\mu_{A,\psi}(\lambda).$$

If $\|\psi\| = 1$ we have $\int_{\sigma(A)} d\mu_{A,\psi}(\lambda) = 1$, so the spectral measure is a probability measure. This makes good on our promise that in quantum theory the spectral theorem would yield a probability measure given a state and an observable - though so far we’ve only shown this for bounded $A$.

The spectral measure also immediately allows us to strengthen the above theorem:

**Spectral Theorem - measurable functional calculus.** Let $A$ be a bounded self-adjoint operator on $H$. Let $\mathcal{B}(\sigma(A))$ denote the bounded Borel functions on $\sigma(A)$ - not identifying those which are equal a.e.. Given $A$ on $H$ and $f \in \mathcal{B}(\sigma(A))$, there is a unique bounded operator $\pi(f) = f(A)$ such that

$$\langle \psi, f(A)\psi \rangle = \int_{\sigma(A)} f(\lambda) d\mu_{A,\psi}(\lambda).$$

Then $\pi : \mathcal{B}(\mathbb{R}) \to \mathcal{L}(H)$ is the unique map satisfying:

1. $\pi$ is a $*$-homomorphism,

2. $\|f(A)\| = \|f\|_\infty$, where the $L^\infty$ norm is defined with the following modified notion of ‘almost everywhere’: we say a property holds a.e.' if for each $\psi \in H$ the property holds except for a set of measure zero relative to $\mu_{A,\psi}$.

3. If $h$ is the function $h(\lambda) = \lambda$, then $h(A) = A$.

Moreover 4. If $f \geq 0$ then $f(A) \geq 0$.

5. $\sigma(f(A)) = f(\sigma(A))$.

6. If $A\psi = \lambda\psi$ then $f(A)\psi = f(\lambda)\psi$.

7. If $AT = TA$ for some $T \in \mathcal{L}(H)$, then $f(A)T = Tf(A)$.

8. Suppose that $f_n \to f$ pointwise a.e.' and $\|f_n\|_\infty$ is bounded. Then $f_n(A) \to f(A)$ strongly.

**Sketch of proof -** For the uniqueness one uses the polarization identity (see below), which implies that a bounded operator is determined by the inner products $\langle \psi, T\psi \rangle$. Property 8 follows from the dominated convergence theorem and the rest of the properties follow
by limiting arguments and the previous version of the spectral theorem, using Lusin’s theorem that given any Borel measure $\mu$ on $\mathbb{R}$, any $f \in B(\mathbb{R})$ is a pointwise $\mu$-a.e. limit of a uniformly bounded sequence of continuous functions. □

Exercise - Prove the polarization identity for any $T \in \mathcal{L}(H)$:

$$\langle \phi, T\psi \rangle = \frac{1}{4} \left[ \{ \langle \psi + \phi, T(\psi + \phi) \rangle - \langle \psi - \phi, T(\psi - \phi) \rangle \} 
- i \{ \langle \psi + i\phi, T(\psi + i\phi) \rangle - \langle \psi - i\phi, T(\psi - i\phi) \rangle \} \right].$$

Exercise - Work out the details of the proof of the above theorem.

The following exercises contain important generalizations of the spectral theorem, and the results will be used in the next section:

Exercise - Prove the following:

**Spectral Theorem for Collections.** Let $A_1, \ldots, A_n$ be a collection of pairwise commuting bounded self-adjoint operators on $H$. Let

$$X = \sigma(A_1) \times \cdots \times \sigma(A_n) \subseteq \mathbb{R}^n.$$  

Let $B(X)$ denote the bounded Borel functions on $X$, again not identifying those which are equal a.e.. There is a map $\pi : B(X) \to \mathcal{L}(H)$ and for each $\psi \in H$ a measure $\mu_\psi$ on $X$ such that

$$\langle \psi, \pi(f)\psi \rangle = \int_X f(x)d\mu_\psi(x)$$

and:

1. $\pi$ is a $*$-homomorphism,

2. $\|\pi(f)\| \leq \|f\|_\infty$, where the $L^{\infty}$ norm is defined with the following modified notion of ‘almost everywhere’: we say a property holds a.e.’ if for each $\psi \in H$ the property holds except for a set of measure zero relative to $\mu_\psi$.

3. If $h : X \to \mathbb{R}$ is the function $h(\lambda_1, \ldots, \lambda_n) = \lambda_i$, then $\pi(h) = A_i$.

4. If $f \geq 0$ then $\pi(f) \geq 0$.

5. If some $T \in \mathcal{L}(H)$ has $A_iT = TA_i$ for all $i$, then $\pi(f)T = T\pi(f)$. 

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6. Suppose that \( f_n \to f \) pointwise a.e. and \( \|f_n\|'_\infty \) is bounded. Then \( f_n(A) \to f(A) \) strongly.

One usually writes \( f(A_1, \ldots, A_n) \) for \( \pi(f) \).

To prove this one can essentially copy the proof for a single operator. However to get the B.L.T. theorem to work don’t use polynomials in the \( A_i \), but instead linear combinations of ‘rectangle functions’. Namely, define a rectangle in \( \mathbb{R}^n \) to be a set of the form \( S_1 \times \cdots \times S_n \) with each \( S_i \) a Borel subset of \( \mathbb{R} \). Define \( \pi \) on characteristic functions of rectangles by:

\[
\pi(\chi_R) = \chi_{S_1}(A_1) \cdots \chi_{S_n}(A_n)
\]

where each \( \chi_{S_i}(A_i) \) is defined by the spectral theorem. Using the spectral theorem for a single bounded operator, and in particular statement 8 of the ‘measurable functional calculus’ version, show that if \( R_i \) is a finite collection of disjoint rectangles then

\[
\| \sum a_i \pi(\chi_{R_i}) \| \leq \| \sum a_i \chi_{R_i} \|_\infty.
\]

Define \( \pi \) on finite linear combinations of rectangle functions by \( \pi(\sum a_i \chi_{R_i}) = \sum a_i \pi(\chi_{R_i}) \), and use the B.L.T. theorem to extend \( \pi \) to \( \mathcal{B}(X) \).

Exercise - A bounded operator is normal if it commutes with its adjoint. E.g., unitary and bounded self-adjoint operators are normal. Show that if \( T \) is normal then \( \text{Re} T = \frac{1}{2}(T + T^*) \) and \( \text{Im} T = \frac{i}{2}(T - T^*) \) are bounded self-adjoint and commute. Use the spectral theorem for collections to prove the following:

**Spectral Theorem for Normal Operators.** Let \( T \) be a normal operator on \( \mathcal{H} \). Let \( \mathcal{B}(\sigma(T)) \) denote the bounded Borel functions on \( \sigma(T) \). There is a map \( \pi : \mathcal{B}(\sigma(T)) \to \mathcal{L}(\mathcal{H}) \) and for each \( \psi \in \mathcal{H} \) a measure \( \mu_{T,\psi} \) on \( X \) such that

\[
\langle \psi, f(T)\psi \rangle = \int_{\sigma(T)} f(\lambda) d\mu_{T,\psi}(\lambda).
\]

and:

1. \( \pi \) is a \( * \)-homomorphism,

2. \( \|f(A)\| = \|f\|'_{\infty} \), where the \( L^\infty' \) norm is defined with the following modified notion of ‘almost everywhere’: we say a property holds a.e. if for each \( \psi \in \mathcal{H} \) the property holds except for a set of measure zero relative to \( \mu_{\psi} \).
3. If \( h \) is the function \( h(\lambda) = \lambda \), then \( h(T) = T \).

4. If \( f \geq 0 \) then \( f(T) \geq 0 \).

5. \( \sigma(f(T)) = f(\sigma(T)) \).

6. If \( T\psi = \lambda \psi \) then \( f(T)\psi = f(\lambda)\psi \).

7. If \( ST = TS \) for some \( S \in \mathcal{L}(H) \), then \( f(S)T = Tf(S) \).

8. Suppose that \( f_n \to f \) pointwise a.e. and \( \|f_n\|_\infty \) is bounded. Then \( f_n(T) \to f(T) \) strongly.

Exercise - Show using the theorem above that if \( T \) is normal, the spectrum of \( T \) is the closure of \( \bigcup_{\psi \in H} \text{supp} \mu_{T,\psi} \). (Hint: show that if \( \lambda \) is not in the closure of \( \bigcup \text{supp} \mu_{T,\psi} \), one can construct \((T - \lambda I)^{-1}\) using the functional calculus.)

Exercise - Use the exercise above to show that the spectrum of a normal operator consists of approximate eigenvalues. Conclude that if \( U \) is unitary, \( \sigma(U) \) is a subset of the unit circle in \( \mathbb{C} \).

Exercise - Assume \( A \in \mathcal{L}(H) \). Use the polarization identity to show that if \( A \geq 0 \), i.e. \( \langle \psi, A\psi \rangle \geq 0 \) for all \( \psi \in H \), then \( A \) is self-adjoint.

Exercise - Show that \( A \geq 0 \) if and only if \( \sigma(A) \) is a subset of \([0, \infty)\) (for one direction use the previous exercise). Show that if \( A \geq 0 \) there exists \( B \in \mathcal{L}(H) \) with \( B \geq 0 \) and \( B^2 = A \).

5. The Spectral Theorem - Unbounded Case

The spectral theorem for unbounded self-adjoint operators proceeds by reduction to the bounded case. We begin by establishing a useful property of self-adjoint operators. We define an operator \( T \) on \( H \) to be symmetric if it is densely defined and for all \( \psi, \phi \in D(T) \) we have

\[
\langle \phi, T\psi \rangle = \langle T\phi, \psi \rangle.
\]

Note that this is equivalent to saying that \( T \subseteq T^* \), since it says that for every \( \phi \in D(T) \) there exists \( \phi' \in H \) (namely \( T\phi \)) such that

\[
\langle \phi, T\psi \rangle = \langle \phi', \psi \rangle,
\]

so that \( D(T) \subseteq D(T^*) \), and \( T^*\phi = \phi' = T\phi \), so that \( T \subseteq T^* \). Thus every self-adjoint operator is symmetric, but as we’ll see the converse does not hold.
Beware: physicists often say self-adjoint when they mean symmetric. This is a result of insufficient education and the following facts: 1) 'sufficiently nice' symmetric operators are self-adjoint, and 2) it's much easier to check if an operator is symmetric than if its self-adjoint - in other words, it can be very difficult to show that a symmetric operator is 'sufficiently nice'.

To check that an operator $T$ is symmetric one must show that $\langle \phi, T\psi \rangle = \langle T\phi, \psi \rangle$ for every $\phi, \psi \in D(T)$; this is often easy to do, and if $T$ is a differential operator it usually amounts to integration by parts. To check that $T$ is self-adjoint one must also show that if $\phi \in H$ has a $\phi' \in H$ such that $\langle \phi, T\psi \rangle = \langle \phi', \psi \rangle$, then $\phi$ lies in $D(T)$. This is usually done indirectly using various theorems, some of which we will discuss later, and many problems in mathematical physics consist of proving that symmetric operators are self-adjoint. This is because the spectral theorem only holds for self-adjoint operators, so self-adjoint operators are much better than merely symmetric operators.

We also note term 'hermitian' is also used by mathematicians to mean symmetric, and by physicists to mean self-adjoint, by which they mean symmetric...we will avoid this term.

We define an operator $T$ to be closed if given $\psi_i \in D(T)$ such that $\psi_i \to \psi \in H$ and $T\psi_i \to \psi'$ then $\psi \in D(T)$ and $T\psi = \psi'$. If $T$ is an operator and $T \subseteq S$ for some operator $S$, we say $S$ is an extension of $T$. If $T$ is an operator with a closed extension we say $T$ is closable. If $T$ is closable the intersection of all closed extensions of $T$ is a closed extension of $T$ called the closure of $T$, written $\overline{T}$.

Exercise - Prove the claim above the intersection of a set of closed extensions is a closed extension.

Exercise - Show that if $T$ is a densely defined operator satisfying $\|T\psi\| \leq C\|\psi\|$ then $T$ has a unique closed extension, and that this extension is a bounded operator.

Exercise - An operator that’s not closable. Let $\psi_n$ be an orthonormal basis for a Hilbert space $H$ and let $\phi \in H$ be a vector that is not a linear combination of finitely many $\psi_n$. Let $D$ be the set of finite linear combinations of $\phi$ and the $\psi_n$, and let $T$ be the operator with domain $D$ defined by:

$$T \left( a\phi + \sum_{n=1}^{N} b_n \psi_n \right) = a\phi.$$  

Show that $T$ is a densely defined linear operator. Show that any closed extension $S$ of $T$ would have $S\phi = \phi$ and $S\phi = 0$, obtaining a contradiction.

Exercise - A nonzero operator whose adjoint is zero (roughly speaking). Let $\psi_n$ be an orthonormal basis for a Hilbert space $H$ and let $D$ be the set of all finite linear combinations of the $\psi_n$. Let $T$ be the operator with domain $D$ defined by:

$$T \left( \sum_{n=1}^{N} a_n \psi_n \right) = (\sum_{n=1}^{N} a_n)\psi_1.$$  

Show that $T$ is a densely defined linear operator. Show that the domain of $T^*$ consists of all vectors orthogonal to $\psi_1$, and that $T^*$ is zero on this domain.
Note that in matrix language

\[ T = \begin{pmatrix} 1 & 1 & 1 & \cdots \\ 0 & 0 & 0 & \cdots \\ 0 & 0 & 0 & \cdots \\ \vdots & \vdots & \vdots & \ddots \end{pmatrix} \]

So why is \( T^* \) zero?

**The Basic Criterion for Self-adjointness.** Let \( T \) be a symmetric operator on a Hilbert space \( H \). Then the following three statements are equivalent:

1. \( T \) is self-adjoint.
2. \( T \) is closed and \( \text{Ker}(T^* \pm i) = \{0\} \).
3. \( \text{Ran}(T \pm i) = H \).

Proof. 1 \( \Rightarrow \) 2: Assume \( T \) is self-adjoint. To show that \( T \) is closed we need to show that if \( \psi_i \to \psi \in H \) and \( T\psi_i \to \psi' \) then \( \psi \in D(T) \) and \( T\psi = \psi' \). Since \( T = T^* \) we only need to show that \( \psi \in D(T^*) \) and \( T^*\psi = \psi' \). In other words, we need to show that for all \( \phi \in D(T) \),

\[ \langle \psi, T\phi \rangle = \langle \psi', \phi \rangle . \]

Since \( T \) is self-adjoint and \( \psi_i \in D(T) \) we have:

\[ \langle \psi, T\phi \rangle = \lim_{i \to \infty} \langle \psi_i, T\phi \rangle = \lim_{i \to \infty} \langle T\psi_i, \phi \rangle = \langle \psi', \phi \rangle \]

as desired.

Next we need to show that if \( T \) is self-adjoint and \( (T + i)\psi = 0 \) or \( (T - i)\psi = 0 \) for some \( \psi \in D(T) \), then \( \psi = 0 \). If \( (T - i)\psi = 0 \),

\[ \langle \psi, i\psi \rangle = \langle \psi, T\psi \rangle = \langle T\psi, \psi \rangle = \langle i\psi, \psi \rangle = -\langle \psi, i\psi \rangle , \]

so \( \psi = 0 \), and the other case works similarly.
2 ⇒ 3: Assume $T$ is symmetric and closed and $\text{Ker}(T^* \pm i) = 0$. We shall show $\text{Ran}(T \pm i)$ is closed and dense in $H$. Since $\text{Ker}(T^* \pm i) = \text{Ran}(T \mp i)^\perp$, and the orthogonal complement of a subspace is empty if and only if the subspace is dense, clearly $\text{Ran}(T \pm i)$ is dense.

To show that $\text{Ran}(T + i)$ is closed we naturally use the hypothesis that $T$ is closed (the proof that $\text{Ran}(T - i)$ is closed is similar). Suppose we have $(T + i)\psi_i \to \phi$ in $H$. We need to show that there exists $\psi \in D(T)$ such that $\phi = (T + i)\psi$. We have

$$0 \leftarrow \|(T + i)\psi_i - (T + i)\psi_j\|^2 = \|T(\psi_i - \psi_j)\|^2 + \|\psi_i - \psi_j\|^2 + \langle i(\psi_i - \psi_j), T(\psi_i - \psi_j) \rangle + \langle T(\psi_i - \psi_j), i(\psi_i - \psi_j) \rangle,$$

and the last two terms cancel because $T$ is symmetric. Thus $\psi_i \to \psi$ and $T\psi_i \to T\psi$ for some $\psi \in D(T)$, since $T$ is closed. Thus we have $\phi = \lim (T + i)\psi_i = (T + i)\psi$.

3 ⇒ 1: Assume that $\text{Ran}(T \pm i) = H$ there exists $\psi \in D(T)$ such that $(T^* - i)\phi = (T - i)\psi$. Since $T \subseteq T^*$ this implies $\phi - \psi \in D(T^*)$ and

$$(T^* - i)(\phi - \psi) = 0.$$  

Since $\text{Ran}(T + i) = H$, $\text{Ker}(T^* - i) = 0$ so $\phi = \psi \in D(T)$, as was to be shown. $\square$

Exercise - Let $X$ be a measure space, $\mu$ a measure on $X$, and $H = L^2(X, \mu)$. Let $f$ be a real-valued measurable function on $X$ and let $M_f$ be the operator with domain equal to $D = \{\psi \in H: f\psi \in H\}$ and for $\psi \in D$ let $M_f\psi = f\psi$. $M_f$ is called a multiplication operator. Show that: 1) $M_f$ is densely defined, 2) $M_f$ is symmetric, and 3) $M_f$ is self-adjoint. One can prove 3) directly from the definitions or using the criterion above.

Spectral Theorem - measurable functional calculus. Suppose that $A$ is a self-adjoint operator on $H$. Then for each $\psi \in H$ there is a spectral measure $\mu_{A,\psi}$, and for each $f \in \mathcal{B}(\mathbb{R})$ there is a unique bounded operator $\pi(f) = f(A)$ such that

$$\langle \psi, f(A)\psi \rangle = \int_{\mathbb{R}} f(\lambda) d\mu_{A,\psi}(\lambda)$$

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and:

1. $\pi$ is a $\ast$-homomorphism,

2. $\|f(A)\| = \|f\|_\infty$, where the $L^\infty$ norm is defined with the modified notion of ‘almost everywhere’ in which we say a property holds a.e.’ if for each $\psi \in H$ the property holds a.e. except for a set of measure zero relative to $\mu_{A,\psi}$.

3. If $f_n \in B(\mathbb{R})$ are functions such that $f_n(\lambda) \to \lambda$ pointwise and $|f_n(\lambda)| \leq |\lambda|$, then $f_n(x)\psi \to A\psi$ for any $\psi \in D(A)$.

4. If $f \geq 0$ then $f(A) \geq 0$.

5. If $A\psi = \lambda\psi$ then $f(A)\psi = f(\lambda)\psi$.

6. Suppose that $f_n \to f$ pointwise a.e.’ and $\|f_n\|_\infty$ is bounded. Then $f_n(A) \to f(A)$ strongly.

Sketch of proof - The idea is to reduce to the bounded case, so we need a function that ‘compresses’ $A$ down to a normal operator without losing any information. For example, the operator consisting of multiplication by $x$ on $L^2(\mathbb{R})$ is self-adjoint and unbounded, and our procedure will ‘compress’ it to the operator consisting of multiplication by $(x + i)(x - i)^{-1}$, which is unitary. Let $\alpha(x) = (x + i)(x - i)^{-1}$. Our idea is to construct $f(A)$ by constructing $(f \circ \alpha^{-1})(\tilde{A})$, where $\tilde{A} = (A + i)(A - i)^{-1}$ will be a well-defined unitary operator, called the Cayley transform of $A$. Composing $f$ with $\alpha^{-1}$ ‘undoes’ the effect of the Cayley transform $A \mapsto \tilde{A}$. We begin with some lemmas:

**Lemma 1.** If $A$ is self-adjoint then $\tilde{A} = (A + i)(A - i)^{-1}$ is a well-defined unitary operator on $H$.

Proof - If $A$ is self-adjoint the maps $(A \pm i): D(A) \to H$ are 1-1 and onto. Thus the composite $(A + i)(A - i)^{-1}$ is a well-defined 1-1 and onto map from $H$ to $H$, and to show that it is unitary it suffices to show that it is norm-preserving.

We need to show that $\|(A + i)(A - i)^{-1}\| = \|\phi\|$ for all $\phi \in H$. We can use the criterion for self-adjointness to write $\phi = (A - i)\psi$ for some $\psi \in D(A)$. Thus it suffices to show that $\|(A + i)\psi\| = \|(A - i)\psi\|$ for all $\psi \in D(A)$. This follows from:

$$\|A \pm i\|\psi\|^2 = \|A\psi\|^2 + \|\psi\|^2,$$

shown earlier in the proof of the criterion for self-adjointness. □

**Lemma 2.** Let $S^1$ denote the unit circle in the complex plane. Given $f \in B(\mathbb{R})$ define $\tilde{f} \in B(S^1)$ by: $\tilde{f}(x) = f(\alpha^{-1}(x))$ on the range of $\alpha$, and $\tilde{f} = 0$ elsewhere. Then the map
\( f \mapsto \tilde{f} \) defines a \(*\)-homomorphism from \( B(\mathbb{R}) \) to \( B(S^1) \) modulo functions that are zero a.e. \( t \), where \( f = 0 \) a.e. means that \( f = 0 \) a.e. with respect to all the measures \( \mu_{\tilde{A},\psi} \).

**Proof** - An easy calculation shows that \( \alpha: \mathbb{R} \to S^1 \), so that \( \tilde{f} \in B(S^1) \). The range of \( \alpha \) is just \( S^1 - \{1\} \). We prove one of the \(*\)-homomorphism properties and leave most of the rest as easy exercises:

\[
(\tilde{f} + \tilde{g})(x) = f(\alpha^{-1}(x)) + g(\alpha^{-1}(x)) = (f + g)(\alpha^{-1}(x)) = (\tilde{f} + \tilde{g})(x).
\]

The only tricky point is that we need to show that \( \tilde{1} \), that is the function equal to 1 on \( S^1 - \{1\} \) and 0 at \( \{1\} \), is equal to 1 a.e. It suffices to show that \( \mu_{\tilde{A},\psi}(1) = 0 \) for all \( \psi \in \mathcal{H} \). We use the:

**Sublemma.** If \( T \) is normal, then \( \lambda \in \sigma(T) \) has nonzero measure with respect to some \( \mu_{T,\psi} \) if and only if \( T\phi = \lambda\phi \) for some nonzero \( \phi \in \mathcal{H} \).

**Proof** - Suppose \( \mu_{T,\psi}(\lambda) \neq 0 \). Let \( \chi_\lambda \) be the characteristic function of \( \{\lambda\} \), and let \( \phi = \chi_\lambda(T)\psi \). We claim that \( \phi \neq 0 \) and \( T\phi = \lambda\phi \). First note that by the spectral theorem

\[
\|\phi\|^2 = \langle \psi, \chi_\lambda(T)^2\psi \rangle \\
= \langle \psi, \chi_\lambda(T)\psi \rangle \\
= \int \chi_\lambda(\lambda')d\mu_{T,\psi}(\lambda) \\
= u_{T,\psi}(\{\lambda\}) > 0,
\]

so \( \phi \neq 0 \). Similarly,

\[
\|(T - \lambda I)\phi\|^2 = \langle \psi, \chi_\lambda(T)^2(T - \lambda)^*(T - \lambda)\psi \rangle \\
= \int \chi_\lambda(\lambda')^2|\lambda' - \lambda|^2d\mu_{T,\psi} \\
= 0,
\]

so \( T\phi = \lambda\phi \).

The converse uses similar ideas and is left as an exercise. \( \Box \)

By this sublemma it suffices to show that there is no nonzero \( \phi \in \mathcal{H} \) with \( \tilde{A}\phi = \phi \). Suppose there was, i.e. suppose that

\[
\phi = (A + i)(A - i)^{-1}\phi.
\]
Let $\psi = (A - i)^{-1}\phi$. Then

$$(A - i)\psi = (A + i)\psi$$

so $\psi = 0$. It follows that $\phi = 0$ because $\text{Ker}(A - i)^{-1} = 0$; this is a contradiction. 

Note that $\sigma(\tilde{A}) \subseteq S^1$ since $\tilde{A}$ is unitary. Thus $\tilde{f}$ restricts to a bounded Borel function on $\sigma(A)$, and given $f \in \mathcal{B}(\mathbb{R})$ we may define $f(A) = \tilde{f}(\tilde{A})$ using the spectral theorem for bounded operators. Most of the statements of the theorem now follow rather simply.

Statement 1 follows from Lemma 2 and the functional calculus for normal operators, e.g.: $(fg)(A) = (\tilde{f}\tilde{g})(\tilde{A}) = \tilde{f}(\tilde{A})\tilde{g}(\tilde{A}) = f(\tilde{A})g(\tilde{A})$.

We define the spectral measures $\mu_{A,\psi}$ using the formula given in the statement of the theorem. Note that if $\mu = \mu_{A,\psi}$ and $\tilde{\mu} = \mu_{\tilde{A},\psi}$ then unraveling the definitions implies

$$\int f d\mu = \langle \psi, f(A)\psi \rangle = \langle \psi, \tilde{f}(\tilde{A})\psi \rangle = \int \tilde{f} d\tilde{\mu}$$

Thus the $\mu$-measure of a set $S$ equals the $\tilde{\mu}$-measure of the set $\alpha(S)$, and statement 2 follows from the corresponding statement for bounded operators.

Statements 4-6 also follow easily from the spectral theorem for bounded operators, and we leave them as exercises. Statement 3 is more touchy since the functions $\{f_n\}$ are not uniformly bounded. We need a lemma:

**Lemma 3.** Let $g(x) = (x + i)^{-1}$. Then $g(A) = (A + i)^{-1}$.

Proof - By definition, $g(A) = \tilde{g}(\tilde{A})$, and it is easily computed that $\tilde{g}(x) = \frac{x - 1}{2ix}$. Thus by the functional calculus for normal operators:

$$g(A) = \tilde{g}(\tilde{A}) = (\tilde{A} - 1)(2i\tilde{A})^{-1}$$

$$= \frac{1}{2i}(1 - \tilde{A}^{-1})$$

$$= \frac{1}{2i}(1 - (A - i)(A + i)^{-1})$$

$$= \frac{1}{2i}((A + i) - (A - i))(A + i)^{-1}$$

$$= (A + i)^{-1};$$

here as always one must take care to justify all the algebraic manipulations involving unbounded operators, but we leave this as an exercise. 

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Now suppose \( f_n \in \mathcal{B}(\mathbb{R}) \) are functions such that \( f_n(\lambda) \to \lambda \) pointwise and \( |f_n(\lambda)| \leq |\lambda| \). Let \( g \) be defined as in the lemma; then the sequence \((f_n + ig) \in \mathcal{B}(\mathbb{R})\) is uniformly bounded and converges to 1 pointwise, so by statement 6 and the lemma we have:

\[
\lim_{n \to \infty} (f_n(A) + i)(A + i)^{-1} \phi = \lim_{n \to \infty} (f_n(A) + i)g(A)\phi = \phi
\]

for all \( \phi \in H \). Given any \( \psi \in D(A) \), we may write \( \psi = (A + i)^{-1} \phi \), and conclude that \( \lim_{n \to \infty} (f_n(A) + i)\psi = (A + i)\psi \), or

\[
\lim_{n \to \infty} f_n(A)\psi = A\psi,
\]

as desired. \( \Box \)

The criterion for self-adjointness is due to von Neumann. The extension of the spectral theorem to unbounded operators is due to von Neumann, ‘Allgemeine Eigenwerttheorie Hermitiescher Functionaloperatoren,’ Math. Ann. 102,49-131 (1929-1930), and Stone, ‘Linear Transformations in Hilbert Spaces and their Applications to Analysis,’ A. M. S. Colloq. Publ. 15 (1932). Carleman in 1923 had pointed out that the spectral theorem would not work for arbitrary symmetric operators.

The spectral theorem has many uses; our first big application of it will be to prove Stone’s theorem in the next section. Note that it fulfills our promise to associate to each state \( \psi \) and observable \( A \) of a quantum system a probability measure \( \mu_{A,\psi} \). We shall now describe some basic consequences of this.

It follows from the theorem above that for all states \( \psi \in D(A) \) we have

\[
\langle \psi, A\psi \rangle = \int \lambda \, d\mu_{A,\psi} (\lambda)
\]

Since \( \mu_{A,\psi} \) represents the probability distribution of values that \( A \) can be measured to have in the state \( \psi \), the quantity \( \langle \psi, A\psi \rangle \) naturally corresponds to the mean value, or expectation value, of the observable \( A \) in the state \( \psi \). Note that this only makes sense for \( \psi \in D(A) \).

Exercise - A projection on a Hilbert space \( H \) is a bounded self-adjoint operator \( P \) with \( P^2 = P \). (Sometimes one considers projections that are not self-adjoint, but in quantum theory one typically assumes self-adjointness, so we build it into the definition.) Show that there is a 1-1 correspondence between projections and closed subspaces of \( H \) given by associating to each projection its range, or conversely to each closed subspace the orthogonal projection onto that subspace. Given a self-adjoint operator \( A \) on \( H \) and a
Borel set $S \subseteq \mathbb{R}$, define the **spectral projection** onto $S$ to be the operator $\chi_S(A)$ given by the functional calculus, where $\chi_S$ is the characteristic function of $S$. Show that the spectral projections are indeed projections and satisfy $\chi_{S \cup T}(A) = \chi_S(A) + \chi_T(A)$ if the sets $S$ and $T$ are disjoint.

Note that $\langle \psi, \chi_S(A)\psi \rangle = \int_S d\mu_{A,\psi}(\lambda)$, so that $\langle \psi, \chi_S(A)\psi \rangle$ corresponds to the probability that the observable $A$ is measured to have a value in the set $S$. Spectral projections are thus very important for computing probabilities.

Exercise - Show that if $\mathbb{R}$ is the disjoint union of a countable collection of Borel sets $S_i$ that the sum of the probabilities $p_i$ that $A$ will be measured to have a value lying in $S_i$ is one.

This is only the beginning of detailed correspondence between concepts from probability theory and from quantum theory. Indeed quantum theory may be regarded as a ‘noncommutative generalization’ of probability theory - the noncommutativity of operator multiplication gives rise to numerous strange effects such as the ‘uncertainty principle’, which we will discuss later.

6. Stone’s Theorem

Stone’s theorem gives a 1-1 correspondence between one-parameter subgroups of unitary operators on $\mathcal{H}$ and self-adjoint operators on $\mathcal{H}$, which in quantum mechanics connects symmetries with their associated conserved quantities. To prove it we need the notion of an ‘essentially self-adjoint’ operator.

Recall that a densely defined operator is closable if it has closed extensions, in which case the closure of the operator is its smallest closed extension. The following exercise shows that symmetric operators are always closable.

Exercise - If $T$ is densely defined then $T^*$ is closed. If $T$ is symmetric it is closable, since $T^*$ is a closed extension of $T$.

A symmetric operator $A$ on a Hilbert space $\mathcal{H}$ is **essentially self-adjoint** if its closure is self-adjoint. The importance of this is that it will be easier to show that operators are es-
sentially self-adjoint than to show that they’re self-adjoint, and for some purposes essential self-adjointness is good enough. One has: self-adjoint ⇒ essentially self-adjoint ⇒ symmetric. There are a lot of useful relationships between these concepts, but for we’ll just mention a couple that we’ll need for Stone’s theorem:

Exercise - If $S$ and $T$ are densely defined and $S \subseteq T$, then $T^* \subseteq S^*$. If $A$ is self-adjoint and $B$ is a symmetric extension of $A$ then $B = A$. (Hint: show $A \subseteq B$ and $B \subseteq B^* \subseteq A^* = A$.)

Exercise - Copy the proof of the basic criterion for self-adjointness to show the following:

**Basic Criterion for Essential Self-Adjointness.** Let $T$ be a symmetric operator on a Hilbert space $H$. Then the following three statements are equivalent:

1. $T$ is essentially self-adjoint.
2. $\ker(T^* \pm i) = \{0\}$.
3. $\text{ran}(T \pm i)$ are dense in $H$.

**Stone’s Theorem.** Let $A$ be a self-adjoint operator and let $U(t) = e^{itA}$. Then $U(t)$ is a strongly continuous one-parameter unitary group, i.e., for all $t$ $U(t)$ is unitary, $U(t+s) = U(t)U(s)$, and $t \to \alpha$ implies $U_{\alpha}(t)\psi \to U(t)\psi$ for all $\psi \in H$. Moreover

$$D(A) = \{ \psi \in H : \lim_{t \to 0} \frac{U(t)\psi - \psi}{t} \text{ exists} \}$$

and for all $\psi \in D(A)$

$$iA\psi = \lim_{t \to 0} \frac{U(t)\psi - \psi}{t}.$$  

Conversely, if $U(t)$ is a strongly continuous one-parameter unitary group, there is a unique self-adjoint operator $A$ such that $U(t) = e^{itA}$. $A$ is called the **infinitesimal generator** of $U(t)$.

Proof - If $A$ is a self-adjoint operator it is an easy exercise using the spectral theorem to show that $e^{itA}$ is a strongly continuous unitary group. If $\psi \in D(A)$ then since $(e^{itA} - 1)/t$
converges pointwise to $\lambda$ as $t \to 0$, with $|(e^{it\lambda} - 1)/t| \leq \lambda$ for $|t| \leq 1$, statement 3 of the spectral theorem for unbounded operators implies that

$$iA\psi = \lim_{t \to 0} \frac{U(t)\psi - \psi}{t}.$$ 

Defining $B$ by

$$D(B) = \{ \psi \in H : \lim_{t \to 0} \frac{U(t)\psi - \psi}{t} \text{ exists} \}$$

and

$$iB\psi = \lim_{t \to 0} \frac{U(t)\psi - \psi}{t},$$

it follows from the above that $A \subseteq B$, and $B$ is symmetric since

$$\langle B\psi, \phi \rangle = \lim_{t \to 0} ([it]^{-1}(U(t) - I)\psi, \phi) = \lim_{t \to 0} (\psi, -(it)^{-1}(U(-t) - I)\phi) = \lim_{t \to 0} (\psi, (it)^{-1}(U(t) - I)\phi) = \langle \psi, B\phi \rangle.$$ 

Since $B$ is a symmetric extension of the self-adjoint operator $A$, the exercise above implies $B = A$. Thus

$$D(A) = \{ \psi \in H : \lim_{t \to 0} \frac{U(t)\psi - \psi}{t} \text{ exists} \}.$$ 

Conversely, suppose that $U(t)$ is a strongly continuous one-parameter unitary group. We will obtain an essentially self-adjoint operator by differentiating $U(t)\psi$ for $\psi$ in a domain of ‘smooth’ vectors, and the closure of this operator will be our self-adjoint $A$.

We use an important trick of Gårding’s to obtain vectors for which $U(t)\psi$ is differentiable. Given $f \in C_0^\infty(\mathbb{R})$ define

$$\phi_f = \int f(t)U(t)\phi \, dt,$$

where the integral can be taken in the sense of Riemann (one can develop a full-blown theory of Lebesgue integration for Hilbert-space-valued functions, but for continuous functions like we have here the Riemann integral will do). Note that

$$\lim_{t \to 0} \left( \frac{U(t) - I}{t} \right) \phi_f = \lim_{t \to 0} \int f(s) \left( \frac{U(t + s) - U(s)}{t} \right) \phi \, ds = \lim_{t \to 0} \int \frac{f(\tau - t) - f(\tau)}{t} U(\tau)\phi \, d\tau = -\int f'(\tau)U(\tau)\phi \, d\tau = \phi_{-f'},$$

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leaving as an exercise the justification for passing the limit through the integral!

We define the Gårding domain by
\[ D = \{ \phi_f : f \in C_0^\infty(\mathbb{R}) \} \]
and define the operator \( B \) on \( D \) by
\[ B\phi = \lim_{t \to 0} \left( \frac{U(t) - I}{it} \right) \phi. \]
We claim: 1) \( D \) is dense, 2) \( B \) is essentially self-adjoint, 3) if \( A = B \) then \( U(t) = e^{itA} \).

For 1), given \( \phi \in H \) choose a sequence \( f_n \in C_0^\infty(\mathbb{R}) \) with \( \text{supp} f_n \subset [-1/n, 1/n] \) and \( \int f_n = 1 \). Then each \( \phi f_n \) is in \( D \) and
\[ \| \phi f_n - \phi \| = \| \int f_n(t)(U(t)\phi - \phi) \, dt \| \leq \sup_{t \in [-1/n, 1/n]} \| U(t)\phi - \phi \|; \]
and since \( U(t) \) is strongly continuous it follows that \( \phi f_n \rightarrow \phi \), so \( D \) is dense.

For 2), first note that \( B \) can be proved symmetric in the same way that the operator \( B \) appearing earlier in the proof was shown to be symmetric. Thus to show that \( B \) is essentially self-adjoint we can use the basic criterion and show that \( \ker(B^* \pm i) = \{ 0 \} \).

We treat \((B^* - i)\), the case of \((B^* + i)\) being analogous. Suppose that \( \psi \in D(B^*) \) has \( B^* \psi = i\psi \). Then for each \( \phi \in D(B) = D \) we have
\[ \frac{d}{dt} \langle U(t)\phi, \psi \rangle = \langle iBU(t)\phi, \psi \rangle = -i(U(t)\phi, B^* \psi) = -i(U(t)\phi, i\psi) = \langle U(t)\phi, \psi \rangle, \]
where in the first line we use the fact that \( U(t)\psi \in D \), so that \( \frac{d}{dt} U(t)\psi = iBU(t)\psi \). Thus the function \( f(t) = \langle U(t)\phi, \psi \rangle \) satisfies \( f' = f \), so \( f(t) = ce^t \). But \( f(t) \) is bounded, since \( U(t) \) is unitary, so we must have \( c = 0 \), and in particular \( \langle \phi, \psi \rangle = 0 \) for all \( \phi \in D \). Since \( D \) is dense this means that \( \psi = 0 \) as desired.

For 3) we compare the two unitary groups \( U(t) \) and \( e^{itA} \). Given \( \phi \in D \) let \( w(t) = U(t)\phi - e^{itA}\phi \). Then \( w(t) \) is differentiable and
\[ w'(t) = iBU(t)\phi - iAe^{itA}\phi = iAU(t)\phi - iAe^{itA}\phi = iAw(t) \]
since \( A|_D = B \). Thus
\[ \frac{d}{dt} |w(t)|^2 = \langle w'(t), w(t) \rangle + \langle w(t), w'(t) \rangle = \langle iAw(t), w(t) \rangle + \langle w(t), iAw(t) \rangle = 0 \]
so \( w(t) = 0 \) for all \( t \), since \( w(0) = 0 \). Thus for all \( \phi \in D \) we have \( U(t)\phi = e^{itA}\phi \). Since \( D \) is dense we must have \( U(t) = e^{itA} \) (two bounded operators which agree on a dense set are equal). □

Note that if \( U(t) \) is a strongly continuous one-parameter unitary group with infinitesimal generator \( A \), for all \( t \) and all \( f \in \mathcal{B}(\mathbb{R}) \) we have

\[
\int f(\lambda)d\mu_{A,\psi}(\lambda) = \langle \psi, f(A)\psi \rangle = \langle e^{itA}\psi, f(A)e^{itA}\psi \rangle = \int f(\lambda)d\mu_{A,U(t)\psi}(\lambda)
\]

so

\[
\mu_{A,U(t)\psi} = \mu_{A,\psi}.
\]

This makes good on our promise that Stone’s theorem would associate to each one-parameter subgroup of symmetries an ‘associated conserved quantity’ - the infinitesimal generator. In the most basic case \( U(t)\psi \) represents the state obtained from \( \psi \) by time translation - that is, waiting \( t \) units of time. Then the infinitesimal generator of \( U(t) \) is called the Hamiltonian, and is often denote \( H \). The fact that the probability distribution of measured values of \( H \) doesn’t change with time, i.e. \( \mu_{H,U(t)\psi} = \mu_{H,\psi} \), is called **conservation of energy**.

Also note that if \( \psi \in D(H) \) then

\[
\frac{d}{dt}U(t)\psi = HU(t)\psi.
\]

This is called (the abstract) **Schrödinger’s equation** and is the basic equation of quantum mechanics, just as Hamilton’s equations are fundamental in classical mechanics. Note also how much Schrödinger’s equation resembles the modern version of Hamilton’s equations:

\[
\frac{d}{dt}U(t)\psi = v_HU(t)\psi.
\]

This is naturally no coincidence; there is a subject of mathematics called ‘geometric quantization’ which arose from investigations into the relationship between these two equations.

Note that one can’t do anything with Schrödinger’s equation until one knows what the Hamiltonian \( H \) is! In what follows we will describe the Hamiltonians for various important physical systems, for example atoms and molecules. We will begin, however, with simpler
examples designed to get a firmer grip on the meaning of the spectral theorem and Stone’s theorem.

7. Self-adjoint Matrices and the Spin 1/2-Particle

So far our only examples of self-adjoint operators are multiplication operators. If $X$ is a measure space, $f$ is a real-valued measurable function on $X$, and $M_f$ is the operator on $L^2(X)$ with domain $D = \{ \psi \in L^2(X) : f\psi \in L^2(X) \}$ defined by:

$$M_f \psi = f\psi, \; \psi \in D$$

then $M_f$ is self-adjoint.

Another sort of example comes from finite-dimensional matrices. If $H$ is finite-dimensional, a dense subspace must be the whole space so a densely defined operator $T$ on $H$ must have $D(T) = H$, and any such operator must be bounded. Thus all self-adjoint operators on finite-dimensional $H$ are bounded. For any self-adjoint operator $A$ on $H$ there is a unitary $U$ such that $UAU^{-1}$ is a diagonal matrix with real coefficients, the eigenvalues of $A$, and $\sigma(A)$ is the set of eigenvalues of $A$, or the set of numbers along the diagonal of $UAU^{-1}$. The following exercise is crucial for getting a feeling for spectral measures:

Exercise - Let $A$ be a self-adjoint operator on $\mathbb{C}^n$ and let $\{ \psi_i \}_{i=1}^n$ be an orthonormal basis of eigenvectors with $A\psi_i = \lambda_i \psi_i$. Let $\psi$ be the state (= unit vector) with $\psi = \sum_{i=1}^n a_i \psi_i$. Show that the spectral measure $\mu_{A,\psi}$ is given by:

$$\mu_{A,\psi} = \sum |a_i|^2 \delta_{\lambda_i}.$$  

The simplest and also most important use of finite-dimensional matrices in quantum theory is in the study of ‘spin-1/2’ particles. During 1921-1924, Stern and Gerlach did experiments in which they sent a beam of silver atoms through an inhomogeneous magnetic field. To their surprise it split into two components. The same effect occurs in a simpler context if one uses electrons rather than silver atoms. Goudsmit and Uhlenbeck hypothesized in 1925 that this strange effect was due to an intrinsic angular momentum of the electron, called ‘spin’. The electrons are either ‘spin up’ or ‘spin down’ (relative to the gradient of the magnetic field). Those with ‘spin up’ are bent one way by the magnetic
field, and those with ‘spin down’ are bent the other. But this happens no matter from which angle the magnetic field is applied! It’s as if one had a spinning tennis ball whose axis of rotation was vertical in any rotated coordinate system, which is clearly impossible.

Wolfgang Pauli, Enrico Fermi, and Paul Dirac developed a mathematical theory of spin-$\frac{1}{2}$ particles (which include protons and neutrons as well as electrons), and we present part of it here in a summarized form. In what follows we ignore the position and momentum aspects of the electron’s state and concentrate on the description of its ‘spin’.

The spin of the electron about any axis is always measured to be $\pm \frac{1}{2} \hbar$, where $\hbar = h/2\pi$ and $h$ is Planck’s constant, a physical constant equal to $6.626196 \times 10^{-27}$ erg sec in the metric system. This constant is so important that mathematical physicists typically work in a modified system of units in which $\hbar = 1$, and we shall do so from now on.

How can the spin of something always be measured to be $\pm \frac{1}{2}$ about any axis? This does not occur in classical mechanics, where the angular momenta (of, say, a spinning tennis ball) about the $x$, $y$, and $z$ axes form a vector $\vec{J} = (J_1, J_2, J_3)$ whose components can be arbitrary real numbers, and transform in the usual way under rotations. We need to specify mathematical descriptions of the states, observables, and symmetry group for the spinning electron. We begin with the states and observables, and discuss the symmetry group later.

In quantum mechanics the states of an electron (ignoring its position and momentum) are described by the unit vectors in a 2-dimensional Hilbert space $\mathcal{H} = \mathbb{C}^2$, and the spin angular momenta about the $x$, $y$, and $z$ axes are observables described by $2 \times 2$ self-adjoint matrices,

$$J_i = \frac{1}{2} \sigma_i,$$

where the $\sigma_i$ are the Pauli matrices:

$$\sigma_1 = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix},$$
$$\sigma_2 = \begin{pmatrix} 0 & -i \\ i & 0 \end{pmatrix},$$
$$\sigma_3 = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}.$$

We develop the properties of these matrices in the following exercises:

Exercise - Show that the each of the matrices $\sigma_i$ has eigenvalues $\pm 1$ (thus the matrices $J_i$ have eigenvalues $\pm \frac{1}{2}$). Determine the eigenvectors of the $\sigma_i$. Note that no two $\sigma_i$ have a (nonzero) eigenvector in common.
Exercise - Show that the operator $J$ defined by

$$J = (J_1^2 + J_2^2 + J_3^2)^{1/2}$$

has $J = \sqrt{3/4}I$. (In general if $\vec{J}$ are the angular momentum operators for a ‘spin-s particle’ one has $J = \sqrt{s(s+1)}I$.)

Exercise - The commutator $[S, T]$ of $S, T \in \mathcal{L}(H)$ is defined to be $ST - TS$. Show that if $A, B \in \mathcal{L}(H)$ are self-adjoint then $i[A, B]$ is self-adjoint. Show that

$$\sigma_i^2 = I$$

and

$$\sigma_i \sigma_j = -\sigma_j \sigma_i = i \sigma_k$$

if $(i, j, k)$ is a cyclic permutation of $(1, 2, 3)$, hence $[J_i, J_j] = i J_k$.

In particular, this implies that real linear combinations of the matrices $i \sigma_1$, $i \sigma_2$, $i \sigma_3$ and the identity $I$ form an algebra over $\mathbb{R}$ that is isomorphic to Hamilton’s quaternions, which recall is an algebra generated by elements $i, j, k$ satisfying

$$i^2 = j^2 = k^2 = -1,$$

$$ij = -ji = k, \; jk = -kj = i, \; ki = -ik = j.$$

In the late 1800’s Hamilton’s quaternions (invented 1843) were in strong competition with Gibbs’ vector notation as a tool for the mathematics of 3-dimensional space. Vector notation, being more elegant for the problems at hand, eventually won out. But it is now clear that quaternions (now often represented as matrices) are necessary for describing certain features of the geometry of $\mathbb{R}^3$. The appropriate generalization to higher dimensions are the ‘Clifford algebras’, known in the four-dimensional case by physicists as ‘Dirac matrices’.

Exercise - Let $\psi \in H$ be the unit vector $(a, b)$, $|a|^2 + |b|^2 = 1$. Determine the spectral measures $\mu_{\sigma_i, \psi}$, using the exercise above on spectral measures.

In particular, the ‘spin-up’ and ‘spin-down’ states

$$\uparrow = \begin{pmatrix} 1 \\ 0 \end{pmatrix}, \; \downarrow = \begin{pmatrix} 0 \\ 1 \end{pmatrix},$$

are the eigenvectors of the operator $J_3$:

$$J_3 \uparrow = \frac{1}{2} \uparrow, \; J_3 \downarrow = -\frac{1}{2} \downarrow,$$
and they correspond to an electron with spin pointing up or down, i.e., rotating counterclockwise or clockwise along the $z$ axis. An electron in any state can be written as a linear combination or, in physics language, superposition of these states:

$$\begin{pmatrix} a \\ b \end{pmatrix} = a\uparrow + b\downarrow ,$$

and has the spectral measure

$$\mu_{J,\psi} = |a|^2 \delta_{-1/2} + |b|^2 \delta_{1/2} .$$

Thus it has probability $|a|^2$ of being measured to have spin $\frac{1}{2}$ and probability $|b|^2$ of being measured to have spin $-\frac{1}{2}$, if one measures its angular momentum along the $z$ axis. If one measures its spin along the $x$ axis, one will measure it to be $\pm \frac{1}{2}$, distributed according to the probability distribution $\mu_{J,\psi}$.

Of course there’s nothing special about the $x, y,$ and $z$ axes. Let $\vec{v} \in \mathbb{R}^3$ be the unit vector $(v_1, v_2, v_3)$. Then the self-adjoint operator corresponding to spin in the $\vec{v}$ direction is

$$\vec{v} \cdot \vec{J} = \sum_{i=1}^{3} v_i J_i .$$

Now we introduce the symmetry group. Recall that the special orthogonal group $SO(n)$ is the group of $n \times n$ real orthogonal matrices with determinant one, and the special unitary group $SU(n)$ is the group of $n \times n$ complex unitary matrices with determinant one. $SO(3)$ is the group of rotations of 3-dimensional space, and is the obvious candidate for being the symmetry group of the spinning electron, but in fact the group of symmetries is $SU(2)$, which we shall show is a double cover of $SO(3)$. That is, there is an onto, two-to-one homomorphism $\alpha : SU(2) \to SO(3)$.

Note first that the group $SU(2)$ has a natural strongly continuous unitary representation $U$ on $H$ given by

$$U(g)\psi = g\psi .$$

(Here as always we give groups of finite-dimensional matrices the topology such that $T_\alpha \to T \iff (T_\alpha)_{ij} \to T_{ij}$ for all $i,j$.) Next note that if $A$ is a self-adjoint operator on $H$ which is traceless, i.e. $\text{tr}A = 0$ (where $\text{tr}A = \sum A_{ii}$), then $e^{itA} \in SU(2)$, by the following:

Exercise - If $H$ is an arbitrary Hilbert space, show that for self-adjoint $A \in \mathcal{L}(H)$,

$$e^{itA} = \sum_{n=0}^{\infty} \frac{(itA)^n}{n!} .$$

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Show that for any self-adjoint $A \in \mathcal{L}(\mathbb{C}^n)$,

$$\det e^{itA} = e^{it \text{tr} A}.$$ 

For any unit vector $\vec{v} \in \mathbb{R}^3$ the one-parameter subgroup $e^{i\theta \vec{v} \cdot \vec{J}}$ ($\theta \in \mathbb{R}$) corresponds to rotations about the $\vec{v}$ axis. In other words, given an electron with state described by $\psi$, if one rotates it by an angle $\theta$ about the $\vec{v}$ axis one obtains an electron with state $e^{i\theta \vec{v} \cdot \vec{J}} \psi$.

Note however that

$$e^{i \theta J_3} = \begin{pmatrix} e^{i \theta/2} & 0 \\ 0 & e^{-i \theta/2} \end{pmatrix},$$

so rotating the state $\psi$ by $2\pi$ gives, not $\psi$, but $-\psi$! Rotating by $4\pi$, we get back to the state $\psi$! This is one reason why it makes sense to call such particles ‘spin-$\frac{1}{2}$’. And it’s a physical fact that if you leave an electron lying on the table and someone sneaks in when you’re not looking and rotates it by 360 degrees, when you come back you can tell - if you have the right apparatus! But if he it turns around twice you can’t.

Here we are glossing over an important subtlety. For every observable $A$ on the Hilbert space $\mathcal{H}$ we have $\mu_A,\psi = \mu_A,-\psi$, so there is no observable property of an electron in isolation that changes upon rotating it 360 degrees. But its interactions with other parts of a larger system may differ.

The spin of a particle can be 0, 1/2, 1, 3/2, etc. Mathematically these various spins correspond to the different irreducible strongly continuous unitary representations of $SU(2)$. Various mesons such as pions, kaons, etc. are spin-zero; electrons, protons, neutrons, neutrinos and various baryons are spin-$1/2$; and photons are spin-one. Assemblages of particles can act like particles with various spins, thus atomic nuclei may be said to have spin $5/2$ and so on.

We have:

**Theorem.** There is an onto, two-to-one homomorphism $\alpha: SU(2) \to SO(3)$ such that the element $e^{i \theta \vec{v} \cdot \vec{J}}$, where $\vec{v}$ is a unit vector, is mapped onto the matrix representing a rotation of angle $\theta$ about the $\vec{v}$ axis.

**Sketch of Proof -** We construct $\alpha$ as follows. There is an isomorphism between $\mathbb{R}^3$ and the $2 \times 2$ self-adjoint traceless matrices given by:

$$\vec{v} \mapsto \vec{v} \cdot \vec{\sigma} = \begin{pmatrix} v_3 & v_1 - iv_2 \\ v_1 + iv_2 & v_2 \end{pmatrix}.$$

Note that

$$||\vec{v}||^2 = -\det(\vec{v} \cdot \vec{J}).$$

The group $SU(2)$ has a representation on the $2 \times 2$ traceless self-adjoint matrices given by

$$(g, A) \mapsto gAg^{-1};$$
we leave it as an exercise to verify this. Via the above 1-1 correspondence this gives rise to a representation of $SU(2)$ on $\mathbb{R}^3$, which we call $\alpha$. Note that
\[
\|V(g)\vec{v}\|^2 = -\det g(\vec{v} \cdot \vec{J})g^{-1} = -\det(\vec{v} \cdot \vec{J}) = \|\vec{v}\|^2,
\]
so that $\alpha(g)$ is an orthogonal $3 \times 3$ matrix.

It’s clear that $\alpha(I) = I$, and $\alpha(gh) = \alpha(g)\alpha(h)$ because
\[
(gh)A(gh)^{-1} = g(hAh)^{-1}g^{-1}.
\]
Thus $\alpha$ is a homomorphism from $SU(2)$ to $O(3)$, the group of $3 \times 3$ orthogonal matrices.

It remains to show that: 1) $\text{Ran}(\alpha) = SO(3)$, 2) $\alpha$ is two-to-one, 3) $\alpha$ maps $e^{i\theta}\vec{v}$ to the rotation of angle $\theta$ about $\vec{v}$.

1) can be shown by a calculation, but it’s easier to use a few facts about these groups. First, since $SU(2)$ is connected and $\alpha$ is continuous, $\alpha$ must map $SU(2)$ into the connected component of $O(3)$ that contains the identity. This is just $SO(3)$; the other connected component consists of matrices with determinant -1. Since every element of $SO(3)$ is a rotation of some angle about some axis, if we show 3) it will follow that $\alpha$ is onto $SO(3)$.

2) is equivalent to saying that the kernel of $\alpha$ consists of two elements. Unraveling the definitions, $g \in SU(2)$ is in $\text{Ker}\alpha$ iff for all traceless self-adjoint matrices $A$ we have $gAg^{-1} = A$. It is easily seen that the only matrices which commute with all traceless self-adjoint matrices are those of the form $cI$. The only ones of these which lie in $SU(2)$ are $\pm I$. Thus $\text{Ker}\alpha = \pm I$.

To show 3) it suffices to show - differentiating with respect to $\theta$ - that:
\[
\frac{d}{d\theta} \alpha(e^{i\theta}\vec{v}, \vec{w})|_{\theta=0} = -\vec{v} \times \vec{w}.
\]
(Here the minus sign is an unimportant consequence of our sign conventions; the main point is that when one rotates the vector $\vec{w}$ about the $\vec{v}$ axis the derivative is $\pm \vec{v} \times \vec{w}$, with the sign depending on whether one rotates clockwise or counterclockwise.) Unraveling the definitions, this is equivalent to:
\[
\frac{d}{d\theta} e^{i\theta\vec{v} \cdot \vec{J}}(\vec{w} \cdot \vec{J})e^{-i\theta\vec{v} \cdot \vec{J}}|_{\theta=0} = -(\vec{v} \times \vec{w}) \cdot \vec{J}.
\]

Doing the differentiation on the left side, it suffices to show:
\[
[i\vec{v} \cdot \vec{J}, \vec{w} \cdot \vec{J}] = -(\vec{v} \times \vec{w}) \cdot \vec{J}.
\]
This follows from the result of an exercise above, that

$$[\sigma_i, \sigma_j] = 2i\sigma_k .$$

\[\square\]

Note also that by general principles the operator $\vec{v} \cdot \vec{J}$ is the conserved quantity corresponding to the one-parameter group $e^{i\theta \vec{v} \cdot \vec{J}}$. This means that the probability distribution of measured spins in the $\vec{v}$ direction doesn’t change if you rotate the electron by a fixed angle about the $\vec{v}$ axis - which makes sense.

Rather than meditating further on the oddity of the quantum theory of spin-$\frac{1}{2}$ particles, let us turn to infinite-dimensional diagonal matrices. Recall that $\ell^2$ is the usual Hilbert space of complex sequences $\psi = (\psi_i)_{i=1}^\infty$ such that $\|\psi\|^2 = \sum |\psi_i|^2$ is finite, with the inner product $\langle \psi, \phi \rangle = \sum \overline{\psi_i}\phi_i$.

**Theorem.** Let $(a_i)_{i=1}^\infty$ be a sequence of real numbers. Let $A$ be the operator on $\ell^2$ with domain

$$D(A) = \{ \psi \in \ell^2 : \sum |a_i\psi_i|^2 < \infty \}$$

given on this domain by

$$(A\psi)_i = a_i\psi_i .$$

Then $A$ is self-adjoint. The restriction of $A$ to the domain $D_0 = \{ \psi \in \ell^2 : \text{only finitely many } \psi_i \text{ are nonzero} \}$ is essentially self-adjoint.

**Proof -** We can think of $\ell^2$ as $L^2(N)$, where $N = 1, 2, 3, \ldots$ with the ‘counting measure’ that assigns to each point the measure 1. Then $A$ is a multiplication operator and is self-adjoint by the general theory of such. To show that $A_0 = A|_{D_0}$ is essentially self-adjoint, we use the basic criterion: it is easily seen that $A_0$ is symmetric, and to show that $\text{Ran}(A \pm i)$ are dense in $H$ it suffices to show that they contain $D_0$, which is dense. Suppose $\psi \in D_0$, and define $\phi$ by $\phi_i = (a_i \pm i)^{-1}\psi_i$. Then $\phi \in D_0$, and $(A_0 \pm i)\phi = \psi$. Thus $\text{Ran}(A \pm i) \supset D_0$. \[\square\]

**Exercise -** Show that if $A$ is the self-adjoint operator on $\ell^2$ given above and $\psi \in \ell^2$ then $\mu_{A,\psi} = \sum |\psi_i|^2\delta_{a_i}$.
The spectrum of an unbounded self-adjoint operator may be defined in more or less
the same way as for bounded operators. We give two definitions and prove that they are
equivalent, after giving an exercise used in the proof:

Exercise - Show that if $T$ is a densely defined operator on $H$ then
$$(T - \lambda I)^* = T^* - \overline{\lambda} I,$$
where $T - \lambda I$ is the operator with domain $D(T)$, defined in the obvious way.

**Theorem.** Let $A$ be a self-adjoint operator on a Hilbert space $H$. Then the following are
equivalent:

1) $\lambda$ is in the closure of $\bigcup \text{supp} \mu_{A,\psi}$.

2) There does not exist $R \in \mathcal{L}(H)$ such that
$$R(A - \lambda I)\psi = \psi \text{ for all } \psi \in D(A) \text{ and } (A - \lambda I)R\psi = \psi \text{ for all } \psi \in H.$$ (Note that for the latter equation to make sense we need to have $\text{Ran}(R) \subseteq D(A)$.)

If either of these holds we say $\lambda$ is in the **spectrum** of $A$, denoted $\sigma(A)$.

**Proof** - 1) $\Rightarrow$ 2): Suppose 1) holds. Then given any $\epsilon > 0$ there exists $\psi \in H$ such
that $\lambda$ is distance $\leq \epsilon$ from $\text{supp} \mu_{A,\psi}$. This implies that the intersection of the interval $I = [\lambda - \epsilon, \lambda + \epsilon]$ with $\text{supp} \mu_{A,\psi}$ is nonzero, so by the spectral theorem
$$\|\chi_I(A)\psi\|^2 = \langle \psi, \chi_I(A)\psi \rangle = \int_I \mu_{a,\psi} > 0.$$ Let $\phi = \chi_I(A)\psi$; by the above $\phi \neq 0$, and clearly $\phi = \chi_I(A)\phi$. It is easily seen using the spectral theorem that $\phi \in D(A)$. Moreover
$$\|(A - \lambda)\phi\|^2 = \langle \phi, (A - \lambda)^2 \chi_I(A)\phi \rangle = \int_I (x - \lambda)^2 d\mu_{A,\phi}(\lambda) \leq \epsilon^2 \int_I d\mu_{A,\phi} = \epsilon^2 \|\phi\|^2.$$ Thus for any $\epsilon > 0$ there exists nonzero $\phi \in D(A)$ such that $\|(A - \lambda I)\phi\| \leq \epsilon\|\phi\|$. It is an easy exercise to show that no bounded $R$ can have $R(A - \lambda I)\phi = \phi$ for all such $\phi$.

2) $\Rightarrow$ 1): If 1) does not hold, then the function $(x - \lambda)^{-1}$ has finite $L^\infty$ norm, and
we can define $R_\lambda$ by the functional calculus to be $(A - \lambda I)^{-1}$. Note that by the functional
calculus we have $R_\lambda^* = R_\lambda^{\ast}$. If we let $f_n$ be a sequence of functions with $f_n(x) \to x$ for
all \( x \in \mathbb{R} \) and and \(|f_n(x)| \leq x\), then by the functional calculus we have \( f_n(A)\psi \to A\psi \) for all \( \psi \in D(A) \), and \( R_\lambda(f_n(A) - \lambda I) \to I \) strongly, hence for all \( \psi \in D(A) \)

\[
R_\lambda(A - \lambda I)\psi = R_\lambda \lim_{n \to \infty} (f_n(A) - \lambda I)\psi = \lim_{n \to \infty} R_\lambda(f_n(A) - \lambda I)\psi = \psi.
\]

Furthermore, if \( \psi \in H \), for all \( \phi \in D(A - \lambda I) \) we have

\[
\langle \psi, \phi \rangle = \langle \psi, R_\lambda(A - \lambda I)\phi \rangle = \langle R_\lambda \psi, (A - \lambda I)\phi \rangle.
\]

hence by the definition of adjoint, \( R_\lambda \psi \in D((A - \lambda I)^*) \), and

\[
(A - \lambda I)R_\lambda \psi = (A - \lambda I)^*R_\lambda \psi = \psi.
\]

Thus 2) does not hold. \( \square \)

We apply this to infinite-dimensional diagonal matrices in the following exercises:

Exercise - Show that a self-adjoint operator \( A \) on \( H \) is bounded if and only if \( \sigma(A) \) is a bounded set, in which case \( \|A\| = r(A) \).

Exercise - If \( A \) is the self-adjoint operator on \( \ell^2 \) given by

\[
(A\psi)_i = a_i \psi_i,
\]

show that \( \sigma(A) = \{a_i\} \). Show that \( A \) is bounded if and only if \( a_i \) is a bounded sequence, in which case \( \|A\| = \sup |a_i| \).

The theorems above have many applications to differential equations. For example:

Exercise - Show that there is an orthonormal basis of \( L^2[0,1] \) consisting of \( \psi_n \in C^\infty[0,1] \) such that

\[
-\frac{d^2}{dx^2} \psi_n = \lambda_n \psi_n
\]

satisfying the Dirichlet boundary conditions:

\[
\psi(0) = \psi(1) = 0.
\]
Let $D$ be the space of finite linear combinations of the $\psi_n$. Show that $-\frac{d^2}{dx^2}$ is essentially self-adjoint on $D$, and let $\Delta_D$ be the closure of this operator. Determine the spectrum of $\Delta_D$.

Exercise - Show that there is an orthonormal basis of $L^2[0,1]$ consisting of $\psi_n \in C^\infty[0,1]$ such that

$$-\frac{d^2}{dx^2}\psi_n = \lambda_n \psi_n$$

satisfying the Neumann boundary conditions:

$$\psi'(0) = \psi'(1) = 0.$$ 

Let $N$ be the space of finite linear combinations of the $\psi_n$. Show that $-\frac{d^2}{dx^2}$ is essentially self-adjoint on $N$, and let $\Delta_N$ be the closure of this operator. Determine the spectrum of $\Delta_N$.

One can then use the functional calculus to define operators such as $e^{-t\Delta}$ (the ‘heat kernel’), which are very useful in mathematics and physics.

8. Translation and Differentiation

Another simple and important self-adjoint operator arises from translation:

**Theorem.** Let $U(t): L^2(\mathbb{R}) \rightarrow L^2(\mathbb{R})$ defined for $t \in \mathbb{R}$ by

$$U(t)\psi(x) = \psi(x + t).$$

Then $U(t)$ is a strongly continuous one-parameter unitary group, and the infinitesimal generator $p$ of $U(t)$ is a self-adjoint operator such that for all $\psi \in D(p)$,

$$p\psi = i^{-1} \frac{d\psi}{dx}$$

where the derivative is taken in the distributional sense.

Proof - Clearly $U(t)$ is a 1-parameter unitary group. Note that to show that a unitary group is strongly continuous it suffices to show that for all $\psi \in H$, $U(t)\psi \rightarrow \psi$ as $t \rightarrow 0$. 

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Given $\psi \in L^2(\mathbb{R})$, for all $\epsilon > 0$ there exists $\phi \in C_0^\infty(\mathbb{R})$ such that $\|\psi - \phi\| \leq \epsilon$. Let $M^2 = \int_{\mathbb{R}} |d\phi/dx|^2$. Then
\[
\|U(t)\psi - \psi\| \leq \|U(t)\psi - U(t)\phi\| + \|U(t)\phi - \phi\| + \|\phi - \psi\|
\leq 2\epsilon + (\int |\phi(x + t) - \phi(x)|^2 dx)^{\frac{1}{2}}
\leq 2\epsilon + |t|M
\]
so if $|t|M \leq \epsilon$ then
\[
\|U(t)\psi - \psi\| \leq 3\epsilon.
\]
Thus $U(t)$ is strongly continuous.

By Stone’s theorem $U(t) = e^{itp}$ for some self-adjoint $p$ and
\[
D(p) = \{ \psi \in L^2(\mathbb{R}) : \lim_{t \to 0} \frac{U(t)\psi - \psi}{t} \text{ exists in } L^2(\mathbb{R}) \}
\]
and
\[
ip = \lim_{t \to 0} \frac{U(t)\psi - \psi}{t}.
\]
Thus we only need to show that if the limit
\[
\lim_{t \to 0} t^{-1}(U(t)\psi - \psi)
\]
exists in $L^2(\mathbb{R})$ then it equals the distributional derivative $d\psi/dx$.

Suppose the limit exists in $L^2(\mathbb{R})$. Then for all test functions $\phi \in C_0^\infty(\mathbb{R})$
\[
\int \frac{d\psi}{dx} \phi dx = -\int \psi \frac{d\phi}{dx} dx
= \lim_{t \to 0} \int t^{-1}(\psi(x + t) - \psi(x)) \phi(x) dx
= \lim_{t \to 0} \int t^{-1}(\overline{\psi(x + t) - \psi(x)}) \phi(x) dx
= \lim_{t \to 0} (t^{-1}(U(t)\psi - \psi), \phi)
= \langle ip\psi, \phi \rangle,
\]
so the distributional derivative $\frac{d\psi}{dx}$ equals $ip\psi$. □

Translation is of course an important physical symmetry, and in the context of physics the infinitesimal generator $p$ is called ‘momentum’. Of course real physics happens in $\mathbb{R}^3$,
not $\mathbb{R}$. Ignoring spin now, the state of a point particle in $\mathbb{R}^3$ is described by a unit vector in the Hilbert space $\mathbf{H} = L^2(\mathbb{R}^3)$, and the group $\mathbb{R}^3$ of translations has a strongly continuous unitary representation on $\mathbf{H}$ given by

$$(U(g)\psi)(x) = \psi(x + g), \ g \in \mathbb{R}^3, \ \psi \in \mathbf{H}.$$ (The proof of the latter is exactly analogous to that of the theorem above.) The infinitesimal generators of translation in the $x$, $y$, and $z$ directions respectively are self-adjoint operators, the $x$, $y$, and $z$ momenta, denoted $p_1$, $p_2$, $p_3$.

We can also obtain the momentum operators, as well as other self-adjoint differential operators, in terms of the Fourier transform. Recall that the Fourier transform is a unitary operator on $L^2(\mathbb{R}^n)$ given by

$$\hat{f}(\vec{k}) = (2\pi)^{-n/2} \int f(\vec{x})e^{-i\vec{k} \cdot \vec{x}}d^n x,$$

where the integral is absolutely convergent for $f \in L^2(\mathbb{R}^n) \cap L^1(\mathbb{R}^n)$ but is defined by a limit for arbitrary $f \in L^2$.

**Theorem.** Let $p$ be the self-adjoint generator of translations on $L^2(\mathbb{R})$. Then the domain of $p$ is equal to $\{f \in L^2(\mathbb{R}): k\hat{f} \in L^2(\mathbb{R})\}$.

Proof - Let $U(t)$ be the one-parameter group of translations on $\mathbf{H} = L^2(\mathbb{R})$, and let $F: \mathbf{H} \to \mathbf{H}$ be the Fourier transform. We have

$$(\widehat{U(t)f})(k) = (2\pi)^{-1/2} \int f(x + t)e^{-ikx}dx = (2\pi)^{-1/2} \int f(x)e^{-ik(x-t)}dx = e^{ikt}\hat{f}(k),$$

or, in short,

$$FU(t) = M_{e^{ikt}}F,$$

or

$$U(t) = F^{-1}M_{e^{ikt}}F,$$

where $M_{e^{ikt}}$ denotes the operator of multiplication by $e^{ikt}$ (which is a unitary operator on $\mathbf{H}$). Thus $\lim_{t \to 0} t^{-1}(U(t) - I)\psi$ exists in $\mathbf{H}$ if and only if $\lim_{t \to 0} t^{-1}(e^{ikt} - I)\hat{\psi}$ exists in $\mathbf{H}$, which is easily seen to hold if and only if $k\hat{\psi} \in \mathbf{H}$. \(\square\)
Exercise - As a spot-check on your real analysis skills, prove the ‘easily seen’ claim above.

Exercise - Determine the spectrum of the operator $p$. Find approximate eigenvectors.

We can use the following theorem to show that $i^{-1} \frac{d}{dx}$ is essentially self-adjoint on $C_0^\infty(\mathbb{R})$. More precisely, if $A$ is a self-adjoint operator on $H$ we say $D \subseteq D(A)$ is a domain of essential self-adjointness for $A$ or a core for $A$ if the closure of $A|_D$ is $A$. (In particular this implies that $A|_D$ is essentially self-adjoint.)

**Theorem.** Let $e^{ita} = U(t)$ be a strongly continuous 1-parameter unitary group on $H$ and suppose $D \subset H$ be a dense space of differentiable vectors for $U(t)$, i.e., if $\psi \in D$ then the limit
\[
\lim_{t \to 0} \frac{U(t)\psi - \psi}{t}
\]
exists in $H$. Suppose also that $D$ is invariant under $U(t)$, i.e., $U(t)D \subseteq D$ for all $t$. Then $D$ is a core for $A$.

**Proof.** The proof is much like part of that of Stone’s theorem. Note that by Stone’s theorem $D \subseteq D(A)$. Next note that $B = A|_D$ is symmetric since $\langle A\psi, \phi \rangle = \langle \psi, A\phi \rangle$ for all $\psi, \phi \in D(A)$, hence for all $\psi, \phi \in D$.

Next we show that $B$ is essentially self-adjoint by the basic criterion, by showing that $Ker(B^* + i) = \{0\}$. We treat $(B^* - i)$, the case of $(B^* + i)$ being analogous. Suppose that $\psi \in D(B^*)$ has $B^*\psi = i\psi$. Then for each $\phi \in D(B) = D$ we have
\[
\frac{d}{dt} \langle U(t)\phi, \psi \rangle = \langle iBU(t)\phi, \psi \rangle
\]
\[
= -i \langle U(t)\phi, B^*\psi \rangle
\]
\[
= -i \langle U(t)\phi, i\psi \rangle
\]
\[
= \langle U(t)\phi, \psi \rangle,
\]
where in the first line we use the fact that $U(t)\psi \in D$, so that $\frac{d}{dt}U(t)\psi = iBU(t)\psi$. Thus the function $f(t) = \langle U(t)\phi, \psi \rangle$ satisfies $f' = f$, so $f(t) = ce^t$. But $f(t)$ is bounded, since $U(t)$ is unitary, so we must have $c = 0$, and in particular $\langle \phi, \psi \rangle = 0$ for all $\phi \in D$. Since $D$ is dense this means that $\psi = 0$ as desired.

Thus we have the following situation: $A$ is self-adjoint, $B \subseteq A$ and $B$ is essentially self-adjoint. We’ll be done if we show the closure of $B$ is $A$. This follows from:

**Lemma.** Suppose $A$ is a self-adjoint operator, $B$ is an essentially self-adjoint operator and $B \subseteq A$. Then $\overline{B} = A$. 40
Proof - Since $A$ is self-adjoint it is closed, and it is easily seen that in general

$$S \subseteq T \implies \overline{S} \subseteq \overline{T},$$

so $\overline{B} \subseteq \overline{A} = A$. To show the opposite inclusion, note that $\overline{B} \subseteq A$ implies $A^* \subseteq \overline{B}^*$, but $A$ and $\overline{B}$ are self-adjoint, so $A \subseteq \overline{B}$. \qed

**Corollary.** The space $C_0^\infty(\mathbb{R})$ is a core for the operator $p$.

Proof - This follows immediately from the theorem. \qed

Differentiation on the unit interval $[0, 1]$ is a more subtle matter. Let $H = L^2[0, 1]$. By Stone’s theorem we know any self-adjoint operator corresponding to $i^{-1} \frac{d}{dx}$ has to generate a one-parameter group of unitaries, and by analogy with the $\mathbb{R}$ case we expect this group to consist of translations, but if one translates a function on $[0, 1]$ it will ‘fall off the edge’!

Suppose one considers, by analogy with the case of $\mathbb{R}$, an operator $B$ by:

$$D(B) = \{ \psi \in H: \psi \in C^\infty[0, 1], \phi(0) = \phi(1) = 0 \}$$

$$B\psi = i^{-1} \frac{d\psi}{dx}.$$ 

Is $B$ essentially self-adjoint? An integration by parts shows that it is symmetric:

$$\langle \psi, B\phi \rangle = \int_0^1 \overline{\psi(x)} i^{-1} \phi'(x) \, dx$$

$$= \int_0^1 i^{-1} \overline{\psi'(x)} \phi(x) \, dx$$

$$= \langle B\psi, \phi \rangle;$$

note that no boundary terms appear in the integration by parts since $\phi(0) = \phi(1) = 0$. Unfortunately $B$ is not essentially self-adjoint because the kernels $\text{Ker}(B^* \pm i)$ are nonzero.

Consider the case of $B^* + i$ (the other is similar). We claim that the function $e^x \in H$ is in $D(B^*)$ and $B^*e^x = -ie^x$, which implies $e^x \in \text{Ker}(B^* + i)$. To prove our claim it suffices to show that

$$\langle e^x, B\phi \rangle = \langle -ie^x, \phi \rangle$$

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for all $\phi \in D(B)$. This follows from integration by parts:

$$
\langle e^x, B\phi \rangle = \int_0^1 e^x i^{-1} \phi'(x) dx \\
= \int_0^1 i^{-1} e^x \phi(x) dx \\
= \langle -ie^x, \phi \rangle.
$$

The moral here is that the domain $D(B)$ is too small, so the function $e^x$ is orthogonal to $(B - i)\phi$ for all $\phi \in H$, giving $B^* + i$ a nonzero kernel.

So why not try a bigger domain? Define $B$ instead by:

$$
D(B) = \{ \psi \in H : \psi \in C^\infty[0, 1] \} \\
B\psi = i^{-1} \frac{d\psi}{dx}.
$$

This domain is unfortunately too big: now $B$ isn’t symmetric! For if $\psi, \phi \in D(B)$, integration by parts gives:

$$
\langle \psi, B\phi \rangle = \int_0^1 \overline{\psi}(x) i^{-1} \phi'(x) dx \\
= \int_0^1 i^{-1} \overline{\psi}(x) \phi(x) dx + \overline{\psi}\phi\big|_0^1 \\
\neq \langle B\psi, \phi \rangle.
$$

Don’t despair, however! Suppose we define $B$ by

$$
D(B) = \{ \psi \in H : \psi \in C^\infty[0, 1] , \psi(1) = \alpha \psi(0) \} , \\
B\psi = i^{-1} \frac{d\psi}{dx},
$$

where $\alpha$ is a complex number with $|\alpha| = 1$. Now $B$ is symmetric because, integrating by parts:

$$
\langle \psi, B\phi \rangle = \int_0^1 \overline{\psi}(x) i^{-1} \phi'(x) dx \\
= \int_0^1 i^{-1} \overline{\psi}(x) \phi(x) dx + \overline{\psi}\phi\big|_0^1 \\
= \langle B\psi, \phi \rangle,
$$

since

$$
\overline{\psi}(1)\phi(1) = |\alpha|^2 \overline{\psi}(0)\phi(0) = \overline{\psi}(0)\phi(0).
$$
Let $\theta$ be any real number with $e^{i\theta} = \alpha$. Note that the functions $\psi_n$ defined by

$$\psi_n(x) = e^{i(2\pi n + \theta)x}$$

are in $D(B)$, and

$$B\psi_n = (2\pi n + \theta)\psi_n.$$  

By Fourier theory it's easy to see that the functions $\psi_n$ form an orthonormal basis of $H$, so finite linear combinations of the $\psi_n$, which lie in $D(B)$, are dense in $H$. In particular, given such a combination, say $\sum a_n \psi_n$, we can write it as

$$\sum a_n \psi_n = (B \pm i) \left( \sum a_n (2\pi n + \theta \pm i)^{-1} \psi_n \right),$$

where the sum on the right side is in $D(B)$, so we see that $\operatorname{Ran}(B \pm i)$ are dense. Thus $B$ is essentially self-adjoint. As one would expect from the above calculations, we have:

**Exercise** - $B$ is a self-adjoint operator with spectrum equal to $\{2\pi n + \theta\}_{-\infty < n < \infty}$.

Thus we see that in this example the expression $i^{-1} \frac{d}{dx}$ may be given many different interpretations as a self-adjoint operator, depending on what domain we choose, and that the spectrum depends on the domain we choose! Note also that the choice of domain is closely related to boundary conditions. This is typical in elliptic operator theory; an elliptic operator on a bounded open set is self-adjoint if we choose the right domain, but its spectrum depends on the domain.

**Exercise** - Describe the unitary group $e^{itB}$ generated by the operator $B$ above, and explain how it depends on $\alpha$. Hint: evaluate $e^{itB} \psi_n$.

9. **The Canonical Commutation Relations**

Ignoring spin, the state of a particle in $\mathbb{R}^3$ is described by a unit vector in $L^2(\mathbb{R}^3)$. The function $\psi$ on $\mathbb{R}^3$ is often called the particle’s wave function. As mentioned above, the momenta are the observables $p_1, p_2, p_3$ generating translations in the $x, y, z$ directions. Other fundamental observables are the position operators, that is, multiplication operators by the coordinate functions $x, y, z$. These self-adjoint operators are denoted $q_1, q_2, q_3$. (The generalization to $\mathbb{R}^n$ should be obvious.)

Recall that $\psi \in L^2(\mathbb{R}^n)$ is said to be in the Schwartz space $\mathcal{S}(\mathbb{R}^n)$ if for all $n, m \geq 0$ we have

$$(1 + r^2)^n (1 + \Delta)^m \psi \in L^2(\mathbb{R}^n).$$
where
\[ \Delta = -\sum_{i=1}^{n} \frac{\partial^2}{\partial x_i^2} \]
is taken in the distributional sense. Another equivalent definition is commonly used:

Exercise - Show that \( \psi \in \mathcal{S}(\mathbb{R}^n) \) if and only if for all \( n, m \geq 0 \) we have
\[
(1 + r^2)^n (1 + \Delta)^m \psi \in L^\infty(\mathbb{R}^n).
\]

Much of the importance of the Schwartz space in operator theory comes from the following facts:

Exercise - Show that if \( \psi \in \mathcal{S}(\mathbb{R}^n) \) then \( p_i \psi \in \mathcal{S}(\mathbb{R}^n) \) and \( q_i \psi \in \mathcal{S}(\mathbb{R}^n) \) for all \( i, 1 \leq i \leq n \).

Exercise - Show that \( \mathcal{S}(\mathbb{R}^n) \) is a core for the operators \( p_i, q_i \). (Hint: use the theorem on cores in the previous section.)

Exercise - Let \( F \) denote the Fourier transform as a unitary operator on \( L^2(\mathbb{R}^n) \). Show that \( F \) maps \( \mathcal{S}(\mathbb{R}^n) \) to itself in a 1-1 and onto manner.

Suppose that \( \psi \in \mathcal{S}(\mathbb{R}^n) \). Then note that for all \( i, j \),
\[
(p_i q_j - q_j p_i)\psi = i^{-1} \left( \frac{\partial}{\partial x_i} x_j \psi - x_j \frac{\partial}{\partial x_i} \psi \right)
= i^{-1} \frac{\partial x_j}{\partial x_i} \psi
= i^{-1} \delta_{ij} \psi.
\]
We also have
\[
(p_i p_j - p_j p_i)\psi = 0 \quad (q_i q_j - q_j q_i)\psi = 0.
\]
These relations, often summarized as
\[
[p_i, q_j] = i^{-1} \delta_{ij} \quad [p_i, p_j] = [q_i, q_j] = 0,
\]
are called the canonical commutation relations. They were discovered by Werner Heisenberg’s supervisor, Max Born, while studying a paper of Heisenberg’s. At the time (1925) Born was one of the few physicists who had a good knowledge of matrices. He interpreted some quantum- mechanical calculations of Heisenberg’s in terms of infinite-dimensional matrices, explaining in this way the odd non-commutativity of multiplication
that Heisenberg had gotten himself into. This made physicists learn about matrices and develop the new ‘matrix mechanics’, as an alternative to Schrödinger’s ‘wave mechanics’. Only later were Heisenberg’s position and momentum operators correlated by Pauli to the operators of multiplication and differentiation by coordinates $x_i$, in which form they had appeared in Schrödinger’s equation for the hydrogen atom (more on which later). Heisenberg had, in effect, been describing differential operators in terms of an orthonormal basis for $L^2(\mathbb{R}^n)$.

10. The Uncertainty Principle

We begin by pointing out how to multiply and add densely defined operators. Given two densely defined operators $A$ and $B$, we define the product $AB$ by:

$$D(AB) = \{ \psi \in D(B) : B\psi \in D(A) \} ,$$

$$AB\psi = A(B\psi) .$$

It is not necessarily the case that $D(AB)$ is dense! This needs to be proved in each case; for example, properties of the Schwartz space make it clear that products of position and momentum operators such as $pq$, $p^2$, $qp^2q$ and so on are all densely defined operators.

Similarly we define the sum $A + B$ by:

$$D(A + B) = D(A) \cap D(B) ,$$

$$(A + B)\psi = A\psi + B\psi .$$

Again it needs to be proved that $D(A + B)$ is dense - two dense subspaces of a Hilbert space can have 0 as their intersection! The properties of the Schwartz space imply that all polynomials in $p$ and $q$ are densely defined operators.

Exercise - Show that if $A$ is a self-adjoint operator then any polynomial in $A$ is densely defined, and any real polynomial in $A$ is essentially self-adjoint.

Exercise - Define $\tilde{q}$ by

$$D(\tilde{q}) = C_0^\infty(\mathbb{R}^n) ,$$

$$(\tilde{q}\psi)(x) = x\psi(x) .$$
Define $\tilde{p}$ by
\[
D(\tilde{p}) = \{ \psi \in L^2(\mathbb{R}) : \hat{\psi} \in C_0^\infty(\mathbb{R}) \}
\]
\[
(\tilde{p}\psi)(x) = i^{-1} \frac{d}{dx}\psi(x).
\]
Show that $\tilde{p}$ and $\tilde{q}$ are essentially self-adjoint and have as closures the usual $p$ and $q$. Show however that
\[
D(\tilde{p}) \cap D(\tilde{q}) = \{0\}
\]
so $\tilde{p} + \tilde{q}$ has as domain only 0! (Hint: for the last part use the fact that if $\hat{\psi} \in L^2(\mathbb{R})$ is compactly supported then $\psi$ is analytic, which follows from the Paley-Wiener theorem.)

Exercise - Show that the operator $\Delta = p_1^2 + p_2^2 + p_3^2$ is essentially self-adjoint. (Hint: show using the Fourier transform that it’s even essentially self-adjoint when restricted to the smaller domain $\mathcal{S}(\mathbb{R}^3)$.)

Given an observable $A$ on $H$ and a state $\psi \in D(A^2)$, we can define the **standard deviation** of $A$ in the state $\psi$, usually written $\Delta A$ even though it depends on $\psi$, by
\[
(\Delta A)^2 = \langle \psi, A^2\psi \rangle - \langle \psi, A\psi \rangle^2.
\]
Note that this corresponds to the usual definition of standard deviation because $\langle \psi, A\psi \rangle$ is the mean value of the observable $A$ in the state $\psi$, as described in section 5.

The **uncertainty principle** gives a lower bound on the product $\Delta A\Delta B$ for non-commuting observables $A$ and $B$:

**Uncertainty Principle.** Let $A$ and $B$ be self-adjoint operators on the Hilbert space $H$, and suppose $\psi \in D(A^2) \cap D(B^2) \cap D(AB) \cap D(BA)$ is a unit vector. Then
\[
\Delta A\Delta B \geq \frac{1}{2} |\langle \psi, [A, B]\psi \rangle|.
\]

Proof - First we note that it suffices to show this for
\[
A' = A - \langle \psi, A\psi \rangle
\]
and
\[
B' = B - \langle \psi, B\psi \rangle
\]
instead. This is because $[A', B'] = [A, B]$ and
\[
\langle \psi, A'\psi \rangle = \langle \psi, A\psi \rangle - \langle \psi, \langle \psi, A\psi \rangle \psi \rangle = 0,
\]
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hence

\[(\Delta A')^2 = \langle \psi, A'^2 \psi \rangle \]
\[= \langle \psi, (A^2 - 2\langle \psi, A\psi \rangle A + \langle \psi, A^2 \psi \rangle) \psi \rangle \]
\[= \langle \psi, A^2 \psi \rangle - \langle \psi, A\psi \rangle^2 \]
\[= (\Delta A)^2 , \]

and similarly for \(B'\).

We have

\[|\langle \psi, [A', B']\psi \rangle| = |\langle A'\psi, B'\psi \rangle - \langle B'\psi, A'\psi \rangle| \]
\[= 2|\text{Im}\langle A'\psi, B'\psi \rangle| \]
\[\leq 2|\langle A'\psi, B'\psi \rangle| \]
\[\leq 2\|A'\psi\|\|B'\psi\| , \]

the last step using Cauchy-Schwarz; but

\[\|A'\psi\| = \langle \psi, A'^2 \psi \rangle^{1/2} = \Delta A' \]

and similarly for \(B'\), so

\[|\langle \psi, [A', B']\psi \rangle| \leq 2\Delta \Delta B , \]

as was to be shown. □

Recall that the momentum and position operators have \([p, q]\psi = i^{-1}\psi\) for all \(\psi \in S(\mathbb{R})\). Starting from this one can show:

Exercise - For all \(\psi \in D(p^2) \cap D(q^2) \cap D(pq) \cap D(qp)\), \([p, q]\psi = i^{-1}\psi .\)

This implies the famous ‘Heisenberg uncertainty principle’:

\[\Delta p \Delta q \geq \frac{1}{2} , \]

for all states \(\psi\) for which both sides make sense. Putting the units back in for once:

\[\Delta p \Delta q \geq \frac{\hbar}{2} . \]

The consequence of this is that one can never measure both the position and momentum of a particle to arbitrary accuracy; accuracy in one must be obtained at the expense of accuracy in the other past a certain point, which is measured by Planck’s constant.
Exercise - Show that the Heisenberg uncertainty principle is sharp: if
\[ \psi(x) = \frac{1}{\sqrt{\pi c}} e^{-x^2/c} \]
show that \[ \|\psi\| = 1 \] and \[ \Delta p \Delta q = \frac{1}{2}. \]

Exercise - Given one of the functions \( \psi \) as in the previous exercise, show that for any \( t \in \mathbb{R} \) the states \( e^{itp}\psi \) and \( e^{itq}\psi \) also have \( \Delta p \Delta q = \frac{1}{2} \). These are translated and ‘boosted’ versions of the state \( \psi \), respectively. Such states are often called ‘wave packets’ since they are states for which the wave function \( \psi \) has minimal uncertainty in position and momentum, hence resembles as closely as possible a classical particle.

11. Quantization, Schrödinger’s Equation, and the Hydrogen Atom

To describe the dynamics of a quantum system, i.e. to specify how states evolve in time, it suffices to specify the Hamiltonian, a self-adjoint operator \( H \) on the Hilbert space \( \mathcal{H} \) whose unit vectors correspond to states: given the state \( \psi \in \mathcal{H} \), \( e^{itH}\psi \) describes the state \( t \) units of time later.

How does one guess the Hamiltonian of a quantum system? I say ‘guess’ because only the real world knows the Hamiltonian. Many quantum systems correspond to classical systems for which one knows a Hamiltonian (see section 2). For example, if one has a classical particle in \( \mathbb{R}^3 \) in a potential \( V: \mathbb{R}^3 \to \mathbb{R} \), the Hamiltonian or energy from classical mechanics is a function \( H: \mathbb{R}^6 \to \mathbb{R} \) given by
\[
H_{\text{class}} = \frac{1}{2m} (p_1^2 + p_2^2 + p_3^2) + V(q_1, q_2, q_3)
\]
where the coordinate functions on the state space \( \mathbb{R}^6 \) are \( (q_1, q_2, q_3, p_1, p_2, p_3) \). The first term in \( H_{\text{class}} \) is called the kinetic energy; \( m \) is a positive number called the mass of the particle. The second term is called the potential energy. An important example is an electron in the electric field of a stationary proton; here the energy is
\[
H_{\text{class}} = \frac{1}{2m_e} (p_1^2 + p_2^2 + p_3^2) - \frac{e^2}{r}
\]
where \( m_e \) is the mass of the electron, \( e \) is the electric charge of the electron (which is minus that of the proton), and \( r = (q_1^2 + q_2^2 + q_3^2)^{1/2} \). A hydrogen atom consists of an electron and
a proton, so classically it would be described by this Hamiltonian. One can show that the solutions of Hamilton’s equations with this Hamiltonian correspond to circular, elliptical, parabolic and hyperbolic motions of the electron about the proton: this is essentially what Newton did in the study of a mathematically isomorphic problem, that of the planets (and comets) moving in the gravitational potential of the sun. (Newton did not, of course, use Hamilton’s formalism!)

As always when speaking of real physical systems, we are glossing over some subtleties. In classical mechanics the earth does not revolve about a stationary sun; rather, the earth and sun revolve about their common center of mass. Since the sun is much more massive than the earth the stationary sun is a good first approximation; similarly, the proton is about 1836 times as massive as the electron. One can solve the problem exactly using a mathematical trick called separation of variables, and it turns out the only effect is to replace the mass of the electron by the reduced mass \( m = \frac{m_e m_p}{m_e + m_p} \), where \( m_p \) is the mass of the proton.

Also, we are temporarily ignoring effects due to the electron’s spin: the proton produces not only an electric field but also a magnetic field, and the latter interacts with the spin of the electron. In general, physical systems are described by a hierarchy of increasingly accurate models which take into account more and more effects, so when someone says ‘the Hamiltonian is...’ this must always be taken with a grain of salt.

One guesses the Hamiltonian for the corresponding quantum system in the most simple-minded manner possible: one formally substitutes the position and momentum operators for the \( q \)’s and \( p \)’s in the expression \( H_{\text{class}} \). Thus for the hydrogen atom the states of the electron are unit vectors in \( \mathcal{H} = L^2(\mathbb{R}^3) \), and the Hamiltonian is

\[
H = \frac{1}{2m} (p_1^2 + p_2^2 + p_3^2) - \frac{e^2}{r},
\]

where in this equation \( p_i \) stands for the momentum operator \( (i^{-1} \partial / \partial x_i) \) and \( r^{-1} \) stands for the operator of multiplication by \( r^{-1} \). Thus \( H \) is a differential operator on \( L^2(\mathbb{R}^3) \). As Schrödinger wrote it, keeping in Planck’s constant (which we normally set equal to 1):

\[
H = -\hbar^2 \frac{\nabla^2}{2m} - \frac{e^2}{r}.
\]

The original Schrödinger equation for the hydrogen atom is:

\[
\frac{1}{i} \frac{d\psi}{dt} = -\frac{\hbar^2}{2m} \nabla^2 \psi - \frac{e^2}{r} \psi,
\]

where \( \psi \in \mathcal{D}(\mathcal{H}) \) and \( \psi(t, \vec{x}) = (e^{itH})\psi(\vec{x}) \).

This process of substituting operator \( p \)’s and \( q \)’s for functions \( p \) and \( q \) is called quantization. Historically it started as a trick with little justification other than that it seemed to work, but now there is a well-developed mathematical theory of it, one aspect of which goes under the name of geometric quantization. (For a mathematically sophisticated
treatment see, for example, Jędrzej Śniatycki, *Geometric Quantization and Quantum Mechanics*.) Note that this process is ambiguous in general, because while $pq =qp$ classically, in quantum mechanics $pq \neq qp$, leading to **operator ordering problems** in the quantization of expressions involving products of $p$’s and $q$’s. These ambiguities can sometimes be resolved through other criteria; note however that in expressions of the form

$$H = \frac{1}{2m}(p_1^2 + p_2^2 + p_3^2) + V(q_1, q_2, q_3)$$

they do not occur! Indeed, one wonders what the history of quantum mechanics would have been like if there has been operator ordering problems in the Hamiltonian of the hydrogen atom (the first system to which quantum theory was successfully applied).

After deciding upon this strategy for quantization, the first problem we must address is: is the quantized Hamiltonian a self-adjoint operator? The next problem is: what is its spectrum? These questions can be answered by computations in the case of the hydrogen atom Hamiltonian, explicitly determining eigenvectors and approximate eigenvectors. For more complicated systems this ‘explicit’ approach is hardly ever possible. So rather than wasting time solving the hydrogen atom, which can be found in any quantum mechanics text, we simply describe the answer.

The answer is: the operator $H_0$ given by

$$D(H_0) = C_0^\infty(\mathbb{R}^3);$$

$$H_0\psi = \left(-\frac{\hbar^2}{2m}\nabla^2 - \frac{e^2}{r}\right)\psi,$$

is essentially self-adjoint. Note that because of the singularity of the potential $e^2/r$ at $r = 0$, some work is needed to show that $H_0$ is really well-defined:

**Exercise** - Show that if $\psi \in D(H_0)$ then $H_0\psi \in L^2(\mathbb{R}^3)$. Show that $H_0$ is symmetric.

Let $\overline{H}_0 = H$. Then the spectrum of the self-adjoint operator $H$ is

$$\{-\frac{me^4}{2\hbar^2n^2}\}_{n=1}^\infty \cup [0, \infty).$$

For $n = 1, 2, 3 \ldots$ there are $n^2$ linearly independent vectors $\psi$ such that

$$H\psi = -\frac{me^4}{2\hbar^2n^2}\psi;$$

the integer $n$ is called the **principal quantum number**. Eigenvectors of the Hamiltonian are also called **bound states**. The $n = 1$ and $n = 2$ states are given as follows:
Exercise - Show that the function

\[ \psi(\vec{x}) = e^{-r/a_0}, \]

where \( r = \|\vec{x}\| \) and \( a_0 = \hbar^2 / me^2 \), satisfies

\[ H\psi = -\frac{me^4}{2\hbar^2}\psi. \]

Thus \( \psi \) is an eigenvector of \( H \) with \( n = 1 \). Show that the functions

\[ \psi(\vec{x}) = xe^{-r/2a_0}, ye^{-r/2a_0}, ze^{-r/2a_0}, \text{ and } (2 - \frac{r}{a_0})e^{-r/2a_0} \]

are eigenvectors of \( H \) with \( n = 2 \). The constant \( a \) is called the **Bohr radius** of the hydrogen atom and is about .529 angstroms, that is, .529 \times 10^{-8} \text{ centimeters.}

In general the states with principal quantum number \( n \) are of the form: \( e^{-r/na_0} \) multiplied by a polynomial of degree \( n \) in \( x, y, z \) and \( r \). Thus the probability of finding the electron in one of these states at a distance greater than \( r \) from the proton decreases rapidly with \( r \). These states correspond physically to states in which the electron is ‘in orbit’ around the proton, so they are also called, especially in chemistry, **orbitals**. The ‘orbit’ idea shouldn’t be taken to seriously, however, because none of the bound states is an eigenvector of the position operators \( q_i \) (which have no eigenvectors!), so the electron does not follow a curve as time goes on: rather, at any time its position is described by the probability density \( |\psi(\vec{x})|^2 \) (this follows from the probability interpretation of the spectral measure). If at time \( t \) the state is an eigenfunction \( \psi \) with \( H\psi = \lambda\psi \), then at time \( t \) the state is given by

\[ \psi(t) = e^{itH}\psi = e^{it\lambda}\psi, \]

so the probability distribution \( |\psi(t, \vec{x})|^2 \) is independent of time! Thus thinking of the electron as a little ball moving around the proton is misleading for such states.

For \( \lambda \geq 0 \) there is no eigenvector \( \psi \in D(H) \) such that \( H\psi = \lambda\psi \), rather, these \( \lambda \) are only approximate eigenvalues of \( H \). The states which are orthogonal to the bound states are called **scattering states** and correspond physically to states in which the electron is ionized, that is, has enough energy to move arbitrarily far from the proton, rather than orbit it. Approximate eigenvectors for \( H \) corresponding to the approximate eigenvalues \( \lambda \geq 0 \) are given as follows:

Exercise - Let \( f \in C_0^\infty(\mathbb{R}) \) be a decreasing function with \( f(r) = 1 \) for \( 0 \leq r \leq 1/2 \) and \( f(r) = 0 \) for \( r \geq 1 \). Given \( \vec{k} \in \mathbb{R}^3 \) let

\[ \psi_n(\vec{x}) = cf(\|\vec{x} - \vec{x}_0\|/n)e^{i\vec{k} \cdot \vec{x}}/\hbar \]
where $\vec{x}_0$ is any point with $\|\vec{x}_0\| = 2n$ and $c$ is chosen such that the $L^2$ norm of $\psi_n$ equals 1. Show that $\psi_n \in D(H)$ and that

$$\lim_{n \to \infty} \|H\psi_n - \lambda \psi_n\| = 0,$$

where $\lambda = \|\vec{k}\|^2/2m$. Conclude that $[0, \infty)$ is in the spectrum of $H$.

12. The Kato-Rellich Theorem

In this section we describe a general method for showing that ‘sufficiently small’ perturbations of self-adjoint operators are self-adjoint. Using it we will show that the hydrogen atom Hamiltonian of the previous section is self-adjoint.

Suppose that $A$ and $B$ are densely defined operators on a Hilbert space $H$ and that:

1) $D(B) \supseteq D(A)$; 2) for some $a, b \in \mathbb{R}$,

$$\|B\phi\| \leq a\|A\phi\| + b\|\phi\|$$

for all $\phi \in D(A)$. Then $B$ is said to be $A$-bounded. The infimum of such $a$ is called the relative bound of $B$ with respect to $A$. Typically $b$ will have to be chosen larger as $a$ is chosen smaller.

We begin with two exercises that we’ll need:

Exercise - Show that if $A$ is a self-adjoint operator then so is $cA$ for any real $c \neq 0$ (where $D(cA) = D(A)$). Conclude that if $A$ is densely defined and symmetric then $A$ is self-adjoint if and only if $\text{Ran}(A \pm is) = H$ for some $s > 0$, in which case it holds for all $s > 0$.

Exercise - Show that if $A$ is densely defined and symmetric then $A$ is essentially self-adjoint if and only if $\text{Ran}(A \pm is)$ is dense in $H$ for some $s > 0$, in which case it holds for all $s > 0$. 

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Kato-Rellich Theorem. Suppose that $A$ is self-adjoint and $B$ is $A$-bounded with relative bound $< 1$. Then $A + B$ is self-adjoint on $D(A)$ and essentially self-adjoint on any core for $A$.

Proof - To show that $A + B$ is self-adjoint we will show that $\text{Ran}(A + B + is) = H$ for all $s \in \mathbb{R}$ with $|s|$ large enough. Since $A$ is self-adjoint, $(A + is): D(A) \to H$ is onto and $(A + B + is)\phi = (I + B(A + is)^{-1})(A + is)\phi$

for all $\phi \in D(A)$. Thus to show that $\text{Ran}(A + B + is) = H$ we need to show that $\text{Ran}(I + B(A + is)^{-1}) = H$. For this it suffices to show that $-1 \notin \sigma(B(A + is)^{-1})$, and since the spectral radius is less than or equal to the norm it suffices to show that $\|B(A + is)^{-1}\| < 1$.

For $\phi \in D(A)$ we have

$$(A + is)\phi = (I + B(A + is)^{-1})(A + is)\phi$$

Thus for $|s|$ sufficiently large, $B(A + is)^{-1}$ has norm less than one, since we can choose $a < 1$.

Suppose $D_0$ is a core for $A$. Then $(A + is)D_0$ is dense in $H$ if $s \neq 0$, $s \in \mathbb{R}$. Choosing $|s|$ large enough, $\text{Ran}(I + B(A + is)^{-1}) = H$ by the above. It follows from the exercise below that $\text{Ran}(A + B + is) = (I + B(A + is)^{-1})(A + is)D_0$ is dense in $H$. Thus $D_0$ is a core for $A + B$.

Exercise - Show that if $T \in \mathcal{L}(H)$, $\text{Ran} T = H$, and $D \subseteq H$ is dense in $H$, then $TD$ is dense in $H$.

To apply the Kato-Rellich theorem to the hydrogen atom we take

$$A = -\frac{1}{2m}\nabla^2, \quad B = -\frac{e^2}{r}.$$ 

The work will consist of showing that $B$ is $A$-bounded. More precisely, recall that $-\nabla^2 = p_1^2 + p_2^2 + p_3^2$ is essentially self-adjoint (an exercise in section 10). We let $\Delta$ denote the closure, a self-adjoint operator on $L^2(\mathbb{R}^3)$. By the following exercise $C^\infty_0$ is a core for $\Delta$.
Exercise - Use the fact that \( C_0^\infty(\mathbb{R}^3) \) is dense in \( S(\mathbb{R}^3) \) to show that \( \Delta \) is essentially self-adjoint when restricted to \( C_0^\infty(\mathbb{R}^3) \). (Hint: express \( \Delta \) in terms of the Fourier transform.)

**Theorem.** Suppose \( V \in L^\infty(\mathbb{R}^3) + L^2(\mathbb{R}^3) \), that is, \( V = V_1 + V_2 \) with \( V_1 \in L^\infty(\mathbb{R}^3) \) and \( V_2 \in L^2(\mathbb{R}^3) \). Then the operator \( \Delta + V \), where here \( V \) denotes the operator of multiplication by \( V \), is essentially self-adjoint on \( C_0^\infty(\mathbb{R}^3) \subset L^2(\mathbb{R}^3) \) and is self-adjoint on \( D(\Delta) \).

**Lemma.** For all \( a > 0 \) there exists \( b \) such that

\[
\| \psi \|_\infty \leq a \| \Delta \psi \|_2 + b \| \psi \|_2
\]

for all \( \psi \in D(\Delta) \subseteq L^2(\mathbb{R}^3) \).

**Proof -** If \( f \) is the Fourier transform of \( \psi \), we have

\[
\| \psi \|_\infty \leq \| f \|_1
\]

by the Riemann-Lebesgue lemma and

\[
\| \psi \|_2 = \| f \|_2 , \| \Delta \psi \|_2 = \| k^2 f \|_2
\]

by the Plancherel theorem, so it suffices to show that for all \( a > 0 \) there exists \( b \) such that

\[
\| f \|_1 \leq a \| k^2 f \|_2 + b \| f \|_2.
\]

Since \( \psi \in D(\Delta) \), \( k^2 f \in L^2(\mathbb{R}^3) \), so for any \( c > 0 \), \( ((ck)^2 + 1) f \in L^2(\mathbb{R}^3) \). It is easily seen that \( ((ck)^2 + 1)^{-1} \in L^2(\mathbb{R}^3) \), so by the Cauchy-Schwartz inequality

\[
\| f \|_1 \leq \| ((ck)^2 + 1)^{-1} \|_2 \| ((ck)^2 + 1) f \|_2
\]

Note that by a change of variables

\[
\| ((ck)^2 + 1)^{-1} \|_2 = \left( \int \frac{d^3 k}{((ck)^2 + 1)^2} \right)^{1/2}
\]

\[
= \left( \int \frac{c^{-3} d^3 k}{(k^2 + 1)^2} \right)^{1/2}
\]

\[
= c^{-3/2} \| (k^2 + 1)^{-1} \|_2 ,
\]

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so letting $d = \|(k^2 + 1)^{-1}\|_2$, we have
\[
\|f\|_1 \leq c^{-3/2}d\|(ck)^2 + 1\|_2 \\
\leq dc^{1/2}\|k^2f\|_2 + dc^{-3/2}\|f\|_2.
\]
Choosing $d$ such that $dc^{1/2} = a$,
\[
\|f\|_1 \leq a\|k^2f\|_2 + b\|f\|_2. \quad \square
\]

Proof of Theorem - It suffices to show that $V$ is $\Delta$-bounded with relative bound less than one. We have
\[
\|V\psi\|_2 = \|(V_1 + V_2)\psi\|_2 \leq \|V_1\|_2\|\psi\|_\infty + \|V_2\|_\infty\|\psi\|_2,
\]
so by the lemma, for any $a > 0$ there exists $b$ such that
\[
\|V\psi\|_2 \leq \|V_1\|_2(a\|\Delta\psi\|_2 + b\|\psi\|_2) + \|V_2\|_\infty\|\psi\|_2,
\]
so $V$ is $\Delta$-bounded with relative bound zero. \quad \square

**Corollary.** The hydrogen atom Hamiltonian,
\[
\frac{1}{2m}\Delta - \frac{e^2}{r},
\]
is essentially self-adjoint on $C^\infty_0(\mathbb{R}^3)$ and self-adjoint on $D(\Delta)$.

Proof - It suffices to show that $\Delta - 2me^2r^{-1}$ has these properties, so by the Kato-Rellich theorem it suffices to show that $r^{-1} \in L^2(\mathbb{R}^3) + L^\infty(\mathbb{R}^3)$. We may write
\[
r^{-1} = \chi_{\{r \leq 1\}}r^{-1} + \chi_{\{r > 1\}}r^{-1},
\]
and the first term is in $L^2$ while the second is in $L^\infty$. \quad \square

**Kato’s theorem** is an extension of the corollary above to Hamiltonians describing arbitrary atoms and molecules, or more generally, collections of charged particles. To describe a system of $n$ particles in $\mathbb{R}^3$ one uses the Hilbert space $L^2(\mathbb{R}^{3n})$, where we give
the coordinates \((x_{1i}, x_{2i}, x_{3i})_{i=1}^{n}\). The position and momentum operators for the \(i\)th particle are
\[
q_{1i} = M_{x_{1i}} , \ q_{2i} = M_{x_{2i}} , \ q_{3i} = M_{x_{3i}}
\]
and
\[
p_{1i} = i^{-1} \frac{\partial}{\partial x_{1i}} , \ p_{2i} = i^{-1} \frac{\partial}{\partial x_{2i}} , \ p_{3i} = i^{-1} \frac{\partial}{\partial x_{3i}} .
\]
If the \(i\)th particle has mass \(m_i\) and charge \(e_i\), the Hamiltonian is
\[
H = \sum_{i=1}^{n} \frac{1}{2m_i} (p_{1i}^2 + p_{2i}^2 + p_{3i}^2) - \sum_{1 \leq i < j \leq n} \frac{e_i e_j}{\| \vec{q}_i - \vec{q}_j \|} .
\]

Exercise - Prove **Kato's Theorem**: the Hamiltonian \(H\) above is essentially self-adjoint on \(C_0^\infty (\mathbb{R}^3_n)\). (Hint: copy the proof for the hydrogen atom).

The kinetic energy term, \(H_0 = \sum_{i=1}^{n} \frac{1}{2m_i} (p_{1i}^2 + p_{2i}^2 + p_{3i}^2)\), is often called the **free Hamiltonian**, since it is the Hamiltonian for \(n\) ‘free’ or non-interacting particles:

Exercise - Let the **classical free Hamiltonian** be given by
\[
H_{class} = \sum_{i=1}^{n} \frac{1}{2m_i} (p_{1i}^2 + p_{2i}^2 + p_{3i}^2) ,
\]
where here the \(p\)'s are functions on \(\mathbb{R}^6_n\) as in section 2. Show that the solutions of Hamilton's equations correspond to particles moving with constant velocity along straight lines.

Showing that the potential energy term or **interaction Hamiltonian**
\[
H_{int} = - \sum_{1 \leq i < j \leq n} \frac{e_i e_j}{\| \vec{q}_i - \vec{q}_j \|}
\]
is \(H_0\)-bounded is a way of making precise the idea that it is in some sense a ‘small perturbation’ of the free Hamiltonian. This situation is typical of mathematical physics: one tries to model a complicated ‘interacting’ system as a perturbation of a simpler ‘free’ system.

This approach has failed so far to prove self-adjointness of the quantum-field-theoretic Hamiltonians describing the basic laws of physics, i.e., quantum electrodynamics, or more comprehensively, the ‘standard model’ of the known particles, interacting via the electromagnetic, weak, and strong forces. These Hamiltonians are not even known to be well-defined operators on a dense domain in some Hilbert space! This is probably **the** problem of mathematical physics today, and the subject of ‘constructive field theory’ is
the hardest testing-ground for new techniques in analysis. There is no reason (any more) to think that in quantum field theory the interaction Hamiltonian is a ‘small’ perturbation of the free Hamiltonian, but there is no other clear way to begin an attack on the problem. Recently, however, I. Segal has proposed a method which amounts to considering the free Hamiltonian as a small perturbation of the interacting Hamiltonian! While this idea is unlikely to crack the problem without a lot of work being put into it, it is an amusing twist.

The Kato-Rellich theorem is originally due to F. Rellich, ‘Störungstheorie der Spektralzerlegung, II’, *Math. Ann.* 116 (1939) 555-570. Kato’s application of this result to atomic Hamiltonians appeared in ‘Fundamental properties of Hamiltonian operators of Schrödinger type’, *Trans. Amer. Math. Soc.* 70 (1951), 195-211. For introductions to mathematically rigorous quantum field theory, try *Introduction to Axiomatic Quantum Field Theory* by Bogolubov, Logunov and Todorov, or *Quantum Physics: A Functional Integral Point of View* by Glimm and Jaffe, or - please forgive the advertising - *Introduction to Algebraic and Constructive Quantum Field Theory* by Baez, Segal and Zhou.
Let $H_{p,c}$ be the following generalization of the hydrogen atom Hamiltonian:

$$H_{p,c} = \Delta + cr^{-p} , \ p \geq 0 .$$

Of course, this is not an honest operator until we also specify its domain. We have seen that for $p = 1$, $H_{p,c}$ is essentially self-adjoint on $C_0^\infty(\mathbb{R}^3)$. We will now study the question: for which $c, p$ is $\Delta + cr^{-p}$ essentially self-adjoint on some reasonable domain? Note that as $p$ increases, the potential $r^{-p}$ becomes more singular at the origin, so our question is an example of the general problem of ‘singular perturbations’: how ’big’ can the interaction Hamiltonian $H_{\text{int}}$ be compared to the free Hamiltonian $H_0$ and still have $H_0 + H_{\text{int}}$ be self-adjoint? Mathematically the situation turns out to be rather complicated, but there are good physical reasons for this complexity.

The important thing is not the particular operator $H_{p,c}$ we’ll be studying but certain general methods that have been developed to study such situations. In this section we consider situations where a symmetric differential operator admits many self-adjoint extensions due to different choices of boundary conditions. We will introduce the ‘deficiency indices’ as a way of measuring how many self-adjoint extensions a symmetric operator has.

First, note the following:

Exercise - Show using the Kato-Rellich theorem that $H_{p,c}$ is essentially self-adjoint on $C_0^\infty(\mathbb{R}^3)$ for $0 \leq p < 3/2$.

Exercise - Show that $r^{-p}\psi$ need not be in $L^2(\mathbb{R}^3)$ for $\psi \in C_0^\infty(\mathbb{R}^3)$ if $p \geq 3/2$.

Thus $H_{p,c}$ is only defined on $C_0^\infty(\mathbb{R}^3)$ for $p \leq 3/2$, in which case it is essentially self-adjoint. One might hope that some other domain would work for larger $p$. In a rather misleading sense this is true for all $p$. First, we note that we can always choose a dense domain on which $H_{p,c}$ is symmetric:

Exercise - Show that the operator $H_{p,c}^0$ defined by

$$D(H_{p,c}^0) = C_0^\infty(\mathbb{R}^3 - 0) ;$$

$$H_{p,c}^0 \psi = (\Delta + cr^{-1})\psi ,$$

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is densely defined and symmetric. (Recall that $C_0^\infty(\mathbb{R}^3 - 0)$ is the space of all compactly supported $C^\infty$ functions on $\mathbb{R}^3$ that vanish in a neighborhood of the origin.)

The next step is to ask: can we find a dense domain on which $H_{p,c}$ is closed and symmetric? The answer is yes, because the closure of a symmetric operator is symmetric:

Exercise - Show that if $T$ is a closable operator, then the graph of the closure of $T$ is the closure of the graph of $T$:  
\[
\{(\psi, T\psi): \psi \in D(T)\} = \overline{\{(\psi, T\psi): \psi \in D(T)\}}.
\]

(Hint: showing $\supseteq$ is easy from the definitions; to show $\subseteq$ show that the right side is the graph of a closed extension of $T$, hence contains the graph of $T$.)

Exercise - Recall the result of an exercise in section 6: every symmetric operator is closable. Show that if $T$ is symmetric then $T$ is symmetric, using the exercise above.

Thus we can take the closure of $H_{p,c}^0$ and get a closed symmetric operator. But is this closure self-adjoint, or if not, does it have self-adjoint extensions? Given a symmetric operator $T$ on $H$, define the deficiency subspaces $H_\pm$ of $T$ by:

\[
H_\pm = \text{Ran}(T \pm i)^\perp = \text{Ker}(T^* \mp i).
\]

Define the deficiency indices of $T$ to be $n_\pm = \dim H_\pm$. The basic criterion for self-adjointness says that a closed symmetric operator $T$ is essentially self-adjoint if and only if $n_+ = n_- = 0$. The following generalization is due to von Neumann:

**Theorem.** Let $T$ be a symmetric operator with deficiency indices $n_+$ and $n_-$. Then $T$ has self-adjoint extensions if and only if $n_+ = n_-$. There is a one-to-one correspondence between self-adjoint extensions of $T$ and unitary maps from $H_+$ to $H_-$.  

**Proof.** We omit the proof; see for example Reed and Simon, Methods of Modern Mathematical Physics, vol. II, Theorem X.2 and corollary thereof. □

We can use this to show that $H_{p,c}^0$ has self-adjoint extensions. Define a conjugation on a Hilbert space $H$ to be a map $K: H \to H$ such that:

\[
K(\psi + \phi) = K\psi + K\phi, \quad K(c\psi) = \overline{c}K(\psi),
\]

where $c \in \mathbb{C}$.
\[ K^2 = I, \|K\phi\| = \|\phi\|, \]
for all \( c \in \mathbb{C} \) and \( \phi, \psi \in \mathcal{H} \). These properties abstract those of the usual complex conjugate.

Exercise - Show that if \( K \) is a conjugation, then for all \( \psi, \phi \in \mathcal{H} \)

\[ \langle K\psi, K\phi \rangle = \langle \psi, \phi \rangle. \]

(hint: use the polarization identity - but \( K \) is conjugate-linear, not linear, so watch out!)

**Theorem.** If \( T \) is a closed symmetric operator on \( \mathcal{H} \) and \( K \) is a conjugation on \( \mathcal{H} \) such that \( K: D(T) \to D(T) \) and \( KT = TK \), then the deficiency indices \( n_{\pm} \) of \( T \) are equal, so \( T \) admits self-adjoint extensions.

Proof - Suppose that \( \psi \in \mathcal{H}_+ \). Then for all \( \phi \in D(T) \),

\[ \langle \psi, (T + i)\phi \rangle = 0, \]

so by the exercise above,

\[ \langle K\psi, (T + i)\phi \rangle = \langle K\psi, (T - i)K\phi \rangle = 0. \]

We thus can conclude that \( K\psi \in \mathcal{H}_- \) if we can show that every vector in \( D(T) \) is of the form \( K\phi \) for some \( \phi \in D(T) \). By our hypotheses \( K: D(T) \to D(T) \), and since \( K^2 = I \), we have \( K^{-1} = K \), so \( K \) is one-to-one and onto from \( D(T) \) to itself. Thus every vector in \( D(T) \) is of the form \( K\phi \) for some \( \phi \in D(T) \).

Thus \( K: \mathcal{H}_+ \to \mathcal{H}_- \). The same argument shows that \( K: \mathcal{H}_- \to \mathcal{H}_+ \). Since \( K = K^{-1} \) this implies that the map \( K: \mathcal{H}_+ \to \mathcal{H}_- \) is 1-1 and onto, so that the dimensions of \( \mathcal{H}_\pm \) are equal, i.e. \( n_+ = n_- \).

**Corollary.** If \( V \) is a real-valued measurable function on \( \mathbb{R}^3 \) such that \( \Delta + V \) is densely defined, then \( \Delta + V \) admits self-adjoint extensions.

Proof - It is easily checked that if \( \Delta + V \) is densely defined then it is symmetric, since \( \Delta \) and \( V \) are. Let \( K: \mathcal{H} \to \mathcal{H} \) be given by:

\[ (K\psi)(x) = \overline{\psi}(x). \]

Then \( K \) is a conjugation, and since \( V \) is real and differentiation commutes with \( K \), \( K \) commutes with \( \Delta + V \). It follows from the theorem above that \( \Delta + V \) admits self-adjoint extensions.
Corollary. The operator $H^0_{p,c}$ admits self-adjoint extensions.

Proof - $H^0_{p,c}$ is symmetric and commutes with $K$ as in the previous corollary. □

One might think that the corollary above solves our problems: can’t one simply take one of the self-adjoint extensions of $H^0_{p,c}$ to be the Hamiltonian? Unfortunately, if there is more than one self-adjoint extension, different choices of self-adjoint extensions generate different one-parameter unitary groups, corresponding to different rules for time evolution. As Reed and Simon put it, ‘different self-adjoint extensions correspond to different physics.’

For example, by analogy with the Hamiltonians we’ve considered so far, the ‘free particle on a half-line’ should be described by the Hilbert space $L^2[0, \infty)$ and the Hamiltonian:

$$-\frac{1}{2m} \frac{d^2}{dx^2}.$$  

But the interpretation of $-d^2/dx^2$ as a self-adjoint operator depends on what domain we choose. If we define $T$ by:

$$D(T) = C_0^\infty(0, \infty);$$

$$T\psi = -\frac{d^2 \psi}{dx^2},$$

clearly $T$ is densely defined, and it is easily seen using integration by parts that $T$ is symmetric:

$$\int_0^\infty \overline{\psi} T \phi = \int_0^\infty \overline{\psi'} \phi' = \int_0^\infty T \psi \phi$$

since $\psi, \phi \in C_0^\infty(0, \infty)$. We can determine the deficiency indices $n_\pm$ of $T$ as follows. If $\psi \in H_+$, then

$$\int_0^\infty \overline{\psi} (T + i) \phi = 0$$

for all $\phi \in D(T) = C_0^\infty(0, \infty)$. Integrating by parts,

$$\int_0^\infty (\psi'' - i \psi') \phi = 0,$$

and since this is true for all $\phi \in C_0^\infty(0, \infty)$, we must have $\psi'' = i \psi$. Thus $\psi \in H_+$ must be a linear combination of the solutions

$$\psi_1(x) = e^{\sqrt{i}x}, \ psi_2(x) = e^{-\sqrt{i}x}.$$

But only $\psi_2 \in L^2(0, \infty)$, since $\psi_1(x)$ grows exponentially as $x \to +\infty$. Therefore $H_+$ is spanned by $\psi_2$, so $n_+ = 1$. Similarly one can show that $H_-$ is spanned by the complex conjugate $\overline{\psi}_2$; this also follows from the theorem on conjugations. Thus $n_- = 1$ as well.
It follows that $T$ has self-adjoint extensions. These are described as follows:

Exercise - Let $T_a, a \in \mathbb{R}$, be the operator given by:

$$D(T_a) = \{ \psi \in C_0^\infty[0, \infty) : \psi'(0) + a\psi(0) = 0 \},$$

$$T_a \psi = -\frac{d^2\psi}{dx^2}.$$

Let $T_\infty$ be given by:

$$D(T_\infty) = \{ \psi \in C_0^\infty[0, \infty) : \psi'(0) = 0 \},$$

$$T_\infty \psi = -\frac{d^2\psi}{dx^2}.$$

Show that $T_a$ and $T_\infty$ are essentially self-adjoint by determining $H_{\pm}$ as above. Conclude that $T_a$ and $T_\infty$ are self-adjoint extensions of $T$.

Thus the choice of self-adjoint extension of $T$ corresponds to a choice of boundary conditions. The choice of boundary conditions determines what happens when a wave function $\psi \in C_0^\infty(0, \infty)$ is evolved in time according to the 1-parameter group generated by the Hamiltonian until it ‘hits’ the ‘wall’ at 0. Dirichlet boundary conditions correspond to the operator $T_\infty$; they produce a situation like that which occurs when one wiggles one end of a rope up and down and the other end is held fixed, so the displacement at that end is zero. If instead the the other end is free to move up and down in a slot, one has Neumann boundary conditions, which correspond to the operator $T_0$. (A real rope satisfies the wave equation, not Schrödinger’s equation, but the similarities rather than the differences concern us here.) If one shakes the rope once, producing a pulse which travels to the other end and bounces back, the pulse will come back upside down given Dirichlet boundary conditions, but right-side up given Neumann boundary conditions!

To verify this, note that $i^{-1} \frac{d}{dx} e^{ikx} = ke^{ikx}$, so that, roughly speaking, $e^{ikx}$ is an eigenvector of $p$ with eigenvalue $k$, or a state with momentum exactly $k$. (Of course $e^{ikx}$ is not in $L^2(0, \infty)$, but one can obtain approximate eigenvectors by ‘cutting it off’ at some large value of $x$.) Show that $\psi = e^{-ikx} + \alpha e^{ikx}$ (suitably cut off) is in the domain of $T_a$ where $\alpha = (ik - a)/(ik + a)$ (where possibly $a = \infty$). Show that $\psi$ is an eigenvector of $T_a$ (ignoring the cutoff). Physically, one interprets $\psi$ as a ‘plane wave’ of momentum $k$ moving to the left superposed with a ‘plane wave’ of momentum $k$ moving to the right, i.e. ‘reflecting’, with the reflected wave having a ‘phase shift’ of $\alpha$. Note that Dirichlet boundary conditions correspond to a phase shift of $-1$ for all momenta, while Neumann boundary conditions correspond to a phase shift of 1 for all momenta.

Something similar happens with $H_{p,c}^0$. If it has more than one self-adjoint extension, different self-adjoint extensions say different things about how the particle ‘pops out’ when it hits the singularity at $r = 0$. In other words, the dynamics is not uniquely determined by $H_{p,c}^0$. The following theorem says when there are such different choices.

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**Theorem.** \( H_{c,p}^0 \) is essentially self-adjoint if \( cr^{-p} \geq \frac{3}{4}r^{-2} \). \( H_{c,p}^0 \) has more than one self-adjoint extension if \( cr^{-p} < \frac{3}{4}r^{-2} \).

**Proof** - The proof is similar to our study of the free particle on the half-line; one obtains a differential equation \( \psi \in H_\pm \) must satisfy and sees which solutions are in \( L^2(\mathbb{R}^3) \). Unfortunately it’s rather lengthy, so we omit it. The proof for \( c \geq 0 \) is Theorem X.11, Reed and Simon, *Methods of Modern Mathematical Physics*, vol. II. The case \( c < 0 \) can be dealt with by also using Theorem XIII.24, Dunford and Schwartz, *Theory of Linear Operators*, vol. II. \( \square \)

Thus we only have essential self-adjointness on \( C_0^\infty(\mathbb{R}^3 - 0) \) if \( p > 2 \) and \( c > 0 \) or \( r = 2 \) and \( c \geq 3/4! \) These conditions are saying, roughly, that the potential \( cr^{-p} \) is sufficiently positive, i.e. repulsive, to make sure that the particle doesn’t blunder into the origin, so that no extra rules are needed to describe how it pops out.

**Exercise** - Let \( T \) be the operator on \( L^2[0,1] \) given by:

\[
D(T) = C_0^\infty(0,1),
\]

\[
T\psi = -\frac{d^2}{dx^2}\psi.
\]

Show that the deficiency indices of \( T \) are \( n_+ = n_- = 2 \) by explicitly determining \( H_\pm \).

(Hint: write down the definition

\[
H_\pm = \{ \psi \in H; \forall \phi \in D(T) \int_0^1 \psi(T \pm i)\phi = 0 \},
\]

and integrate by parts to get a differential equation \( \psi \) must satisfy.)

**Exercise** - Let \( T \) be the operator on \( L^2[0,\infty) \) given by:

\[
D(T) = C_0^\infty(0,\infty),
\]

\[
T\psi = i^{-1}\frac{d\psi}{dx}.
\]

Show that the deficiency indices of \( T \) are \( n_+ = 1, n_- = 0 \). Show that \( T \) has no self-adjoint extensions. Show that \( T \) is a closed symmetric operator that is not self-adjoint!
14. Singular Perturbations and Positivity

Recall that an operator $T$ is bounded below if for some $c$, $T \geq cI$, that is, for all $\psi \in D(T)$
\[ \langle \psi, T\psi \rangle \geq c\|\psi\|^2. \]

The operator $T$ is nonnegative if $T \geq 0$, that is,
\[ \langle \psi, T\psi \rangle \geq 0 \]
for all $\psi \in D(T)$. The operator $T$ is positive if $\langle \psi, T\psi \rangle > 0$ for all nonzero $\psi \in D(T)$. In physics, ‘well-behaved’ Hamiltonians are bounded below; roughly speaking, systems with Hamiltonians not bounded below would be ‘unstable’, being able to emit an infinite amount of energy. This idea is made precise in thermodynamics, where it is seen that at low temperatures a system in thermal equilibrium with its surroundings is in the state of least possible energy (roughly), and that ‘equilibrium’ is only well-defined for Hamiltonians that are bounded below. Note that the free Hamiltonian for a particle in $\mathbb{R}^3$, $\frac{1}{2m}\Delta$, is bounded below, in fact positive, for $m > 0$, since
\[ \langle \psi, \Delta\psi \rangle = \int |\nabla \psi|^2 = \int k^2 |\hat{\psi}|^2 d^3k > 0 \]
for all nonzero $\psi \in D(\Delta)$. (Here the gradient is taken in the distributional sense.) This is one reason why Schrödinger’s equation is only a reasonable description of particles with positive mass.

Photons and neutrinos are thought to have zero mass, although present experimental evidence does not rule out a small mass for neutrinos. Any reasonable description of massless particles require special relativity to be taken into account; such particle always move at the speed of light, and at speeds near that of light the classical Hamiltonian is seen to be only an approximation to a ‘relativistic’ Hamiltonian. Electrons in heavy atoms move at speeds near enough to that of light for these corrections to become important; in the hydrogen atom relativistic effects are quite small (the speed of the electron is about 1/137 times that of light) but very well measured, and provide important confirmation of quantum electrodynamics, the theory of electromagnetism which takes quantum theory and special relativity into account.

Mathematically, bounded below operators are nice in many ways. First of all, it follows from polarization that if $T$ is positive then $T$ is symmetric. As we will show in the next section, there is a nice self-adjoint extension of any positive operator, called the Friedrichs extension.

We will apply these ideas to our Hamiltonian
\[ D(H^0_{p,c}) = C^\infty_0(\mathbb{R}^3 - 0), \]
\[ H_{p,c}^0 \psi = (\Delta + cr^{-p})\psi \, . \]

**Theorem.** \( H_{p,c}^0 \) is not bounded below if \( c < 0 \) and \( p > 2 \).

Proof - Let \( \psi \in C_0^\infty (\mathbb{R}^3 - 0) \) be nonzero. Then clearly \( \langle \psi, \Delta \psi \rangle > 0 \) while if \( c < 0 \), \( \langle \psi, cr^{-p}\psi \rangle < 0 \). Thus for some \( \epsilon > 0 \),

\[ \langle \psi, \Delta \psi \rangle = \epsilon \| \phi \|^2 \]

and for some \( \delta > 0 \),

\[ \langle \psi, \psi \rangle = -\delta \| \phi \|^2 \]

Now let \( \psi_a \) be the ‘dilation’ of \( \psi \) given by

\[ \psi_a(x) = \psi(ax) \, . \]

Note that by a change of variables

\[ \| \psi_a \|^2 = \int |\psi(ax)|^2 d^3x = \int |\psi(y)|^2 a^{-3} d^3y = a^{-3} \| \psi \|^2 \, . \]

On the other hand,

\[ \langle \psi_a, \Delta \psi_a \rangle = \int \bar{\psi_a}(\Delta \psi_a) \]

\[ = a^2 \int \bar{\psi}(ax)(\Delta \psi)(ax) d^3x \]

\[ = a^2 \int \bar{\psi}(y) \Delta \psi(y) a^{-3} d^3y \]

\[ = a^{-1} \langle \psi, \Delta \psi \rangle \]

\[ = a^{-1} \epsilon \| \psi \|^2 \, , \]

while

\[ \langle \psi_a, cr^{-p}\psi_a \rangle = \int c \| x \|^{-p} \bar{\psi}(ax)\psi(ax) d^3x \]

\[ = \int c \| y/a \|^{-p} \bar{\psi}(y)\psi(y) a^{-3} d^3y \]

\[ = a^{p-3} \langle \psi, cr^{-p}\psi \rangle \]

\[ = -a^{p-3} \delta \| \psi \|^2 \, . \]

It follows that for any \( a > 0 \), \( \psi_a \) satisfies

\[ \langle \psi_a, H_{p,c}^0 \psi_a \rangle = \langle \psi_a, (\Delta + cr^{-p})\psi_a \rangle \]

\[ = (a^{-1} \epsilon - a^{p-3} \delta) \| \psi \|^2 \]

\[ = (a^2 \epsilon - a^p \delta) \| \psi_a \|^2 \, . \]
It follows that if \( p > 2 \), no matter how large \( M \) is we can choose \( a \) large enough that

\[
\langle \psi_a, H^0_{p,c} \psi_a \rangle \leq -M \| \psi_a \|^2.
\]

Thus \( H^0_{p,c} \) cannot be bounded below if \( c < 0 \) and \( p > 2 \).  

This is a typical example of a scaling argument: many counterexamples can be produced by considering the family of all dilations \( \psi_a \) of a given function. Here the proof worked by showing that for functions concentrated sufficiently near the origin, namely \( \psi_a \) with \( a \) large, \( cr^{-p} \psi_a \) is bigger than \( \Delta \psi_a \) if \( p > 2 \). The exponent \( p = 2 \) is an example of a critical exponent, at which the proof breaks down because \( \Delta \) and \( r^{-2} \) have the same behavior under scaling. In physics terminology, \( \Delta \) and \( r^{-2} \) both have dimensions of length\(^{-2} \). Many subtle and interesting problems in mathematics and physics arise from critical exponents.

**Theorem.** \( H^0_{p,c} \) is bounded below if \( c > 0 \) or \( p < 2 \). If \( p = 2 \), \( H^0_{p,c} \) is nonnegative if \( c \geq -1/4 \), and unbounded below if \( c < 1/4 \).

**Proof -** \( H^0_{p,c} \) is clearly positive if \( c > 0 \), since \( \Delta \) is positive and \( cr^{-p} \) is nonnegative. For the rest, we use the following beautiful inequality:

**Lemma.** Let \( \psi \in C_0^\infty(\mathbb{R}^3) \). Then

\[
\int_{\mathbb{R}^3} \frac{1}{4r^2} |\psi(x)|^2 \, d^3x \leq \int_{\mathbb{R}^3} |\nabla \psi(x)|^2 \, d^3x.
\]

**Proof -** We may suppose without loss of generality that \( \psi \) is real-valued. Note that if \( r > 0 \)

\[
\nabla (r^{1/2} \psi) = r^{1/2} \nabla \psi + \frac{1}{2} r^{-3/2} \vec{x} \psi
\]

so

\[
|\nabla \psi|^2 = (r^{-1/2} \nabla (r^{1/2} \psi) - \frac{1}{2} r^{-2} \vec{x} \psi)^2
\]

\[
\geq -r^{-5/2} \vec{x} \cdot \nabla (r^{1/2} \psi) + \frac{1}{4} r^{-2} |\psi|^2,
\]

where we are using the fact that \((\vec{v} - \vec{w})^2 \geq \vec{v}^2 - 2\vec{v} \cdot \vec{w}\). Recalling that \( \vec{x} \cdot \nabla f = r \frac{\partial f}{\partial r}, \) we obtain

\[
|\nabla \psi|^2 \geq -r^{-3/2} \psi \frac{\partial}{\partial r} (r^{1/2} \psi) + \frac{1}{4} r^{-2} |\psi|^2.
\]

and using the fact that

\[
r^{-3/2} \psi \frac{\partial}{\partial r} (r^{1/2} \psi) = r^{-1} \frac{\partial \psi}{\partial r} + \frac{1}{2} r^{-2} \psi^2
\]

\[
= \frac{1}{2} r^{-2} \frac{\partial}{\partial r} (r \psi^2)
\]

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we obtain

$$|\nabla \psi|^2 \geq -\frac{1}{2r^2} \frac{\partial}{\partial r} (r|\psi|^2) + \frac{1}{4r^2} |\psi|^2.$$  

Integrating, this yields

$$\int |\nabla \psi|^2 \geq \int \frac{1}{4r^2} |\psi|^2 - \frac{1}{2} \int \frac{\partial}{\partial r} (r|\psi|^2) \sin^2 \theta \, drd\theta d\phi ,$$

and doing the second integral by parts we get zero, so the lemma is proved. \(\square\)

Proof of Theorem, continued - If \(p = 2\), the lemma above immediately implies that \(H_{p,c}^0\) is nonnegative if \(c \geq -1/4\). We treat the case \(c < 1/4\) in the exercise below. We treat the case \(p < 2, c < 0\) as follows. Given \(a > 0\), choose \(\epsilon\) such that \(r^{-p} \leq ar^{-2}\) for all \(r \leq \epsilon\). Then if \(\psi \in C_0^\infty(\mathbb{R}^3)\), we use the lemma to show that:

$$\int_{\mathbb{R}^3} r^{-p} |\psi|^2 \leq \int_{r \leq \epsilon} r^{-p} |\psi|^2 + \int_{r > \epsilon} r^{-p} |\psi|^2 \leq a \int_{\mathbb{R}^3} r^{-2} |\psi|^2 + \epsilon^{-p} \int_{\mathbb{R}^3} |\psi|^2 \leq 4a \int_{\mathbb{R}^3} |\nabla \psi|^2 + \epsilon^{-p} \int_{\mathbb{R}^3} |\psi|^2$$

In particular, if \(\psi \in C_0^\infty(\mathbb{R}^3 - 0)\), we have

$$\langle \psi, H_{p,c}^0 \psi \rangle = \int cr^{-p} |\psi|^2 + |\nabla \psi|^2 \geq (1 + 4ac) \int |\nabla \psi|^2 + \epsilon^{-p} \int |\psi|^2 ,$$

so if we choose \(a\) small enough that \(1 + 4ac \geq 0\),

$$\langle \psi, H_{p,c}^0 \psi \rangle \geq \epsilon^{-p} \int |\psi|^2 = \epsilon^{-p} \|\psi\|^2 ,$$

which shows that \(H_{p,c}^0\) is bounded below. \(\square\)

Exercise - Show that if \(p = 2\) and \(c < -1/4\) then \(H_{p,c}^0\) is not bounded below, as follows.

Assuming \(0 < a < 1/2\), show that the function \(\psi(r) = r^{-a}\) satisfies

$$|\nabla \psi| = \frac{a}{r} |\psi| ,$$

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and the integrals
\[ \int_{\mathbb{R}^3} |\nabla \psi|^2, \int_{\mathbb{R}^3} \frac{1}{r^2} |\psi|^2 \]
converge and satisfy
\[ \int |\nabla \psi|^2 = \int \frac{a^2}{r^2} |\psi|^2. \]
Show that for any \( \epsilon > 0 \) there exists \( \phi \in C_0^\infty(\mathbb{R}^3 - 0) \), obtained by making \( \psi \) smooth and zero near the origin, such that
\[ \int |\nabla \phi|^2 = \int (\frac{a^2}{r^2} + \epsilon) |\phi|^2. \]
Show that if \( c < -1/4 \), for any \( M > 0 \) we can choose \( a \) near 1/2 and \( \epsilon \) small, so that
\[ \langle \phi, H_{p,c}^0 \phi \rangle < -M \| \phi \|^2. \]

15. Singular Perturbations and Sesquilinear Forms

Fundamental problems in mathematical physics often require for their solution the generalization of existing concepts. Often these generalizations arise from the need to deal with more singular objects than the existing framework is capable of handling. For example, distributions arose as a generalization of functions in order to make sense of objects such as the Dirac delta ‘function’. Densely defined operators arose as a generalization of bounded operators because most of the Hamiltonians arising in physics are not bounded. It turns out that there is a very useful further generalization of densely defined operators, called ‘sesquilinear forms’, which are useful in dealing with singular perturbations of self-adjoint operators. In this section we will develop a bit of the theory of sesquilinear forms and apply it to the Hamiltonian \( H_{p,c} \) we have been discussing.

A \textbf{sesquilinear form} on a Hilbert space \( \mathbf{H} \) is a map \( q: Q(q) \times Q(q) \to \mathbb{C} \), linear in the second argument and conjugate-linear in the first, where \( Q(q) \) is a dense subspace of \( \mathbf{H} \), called the \textbf{form domain} of \( q \). (‘Sesqui’ is Greek for ‘one and a half’.)

Given a densely defined operator \( T \) there is a sesquilinear form \( q \) with \( Q(q) = D(T) \) given by:
\[ q(\psi, \phi) = \langle \psi, T\phi \rangle. \]
Not all sesquilinear forms arise in this way, however; a different sort of example is the sesquilinear form \( q \) with \( Q(q) = C_0^\infty(\mathbb{R}) \subset L^2(\mathbb{R}) \) given by

\[
q(\psi, \phi) = \overline{\psi}(0)\phi(0).
\]

Roughly speaking, \( q \) corresponds to ‘multiplication by the delta function at zero’, since

\[
q(\psi, \phi) = \int \overline{\psi}(x)\delta(x)\phi(x)dx,
\]

though of course there is no such multiplication operator, since \( \delta \) is a distribution but not a measurable function. Thus sesquilinear forms are able to describe more singular situations than operators can.

The product of a sesquilinear form \( q \) by a complex number \( c \) is defined to have form domain \( Q(cq) = Q(q) \) and on that domain is given by:

\[
(cq)(\psi, \phi) = cq(\psi, \phi).
\]

The adjoint of a sesquilinear form \( q \) is defined to have form domain \( Q(q^*) = Q(q) \) and on that domain is given by:

\[
q^*(\psi, \phi) = q(\phi, \psi).
\]

Note that we always have \( q^{**} = q \), contrary to the situation for densely defined operators, which is a relief. The sum of two sesquilinear forms \( q \) and \( q' \) is defined to have form domain \( Q(q + q') = Q(q) \cap Q(q') \) and on that domain is given by:

\[
(q + q')(\psi, \phi) = q(\psi, \phi) + q'(\psi, \phi).
\]

The sesquilinear form \( I \) is defined by \( Q(I) = H \) and

\[
I(\psi, \phi) = \langle \psi, \phi \rangle.
\]

One problem with sesquilinear forms is that there is no good general notion of a product of two forms, just as there is no general way to multiply distributions; in certain circumstances, however, a useful product exists. See John Baez, ‘Wick Products of the Free Bose Field’, to appear in Jour. Funct. Anal. For an detailed introduction to sesquilinear forms (also sometimes called quadratic forms) see Barry Simon, Quantum Mechanics for Hamiltonians Defined as Quadratic Forms.

A very important fact is that if \( A \) is an unbounded self-adjoint operator and \( q \) is the sesquilinear form with \( Q(q) = D(A) \) and \( q(\psi, \phi) = \langle \psi, A\phi \rangle \), the form \( q \) can be extended to a larger domain.
**Theorem.** Let $A$ be a self-adjoint operator on a Hilbert space $H$. There is a unique sesquilinear form $q_A$ such that

$$Q(q_A) = \{ \psi \in H : \int_{\sigma(A)} |\lambda| d\mu_{A,\psi}(\lambda) < \infty \}$$

and

$$q_A(\psi, \psi) = \int_{\sigma(A)} \lambda d\mu_{A,\psi}(\lambda)$$

for all $\psi \in Q(q_A)$. $D(A) \subseteq Q(q_A)$ and $q_A(\psi, \phi) = \langle \psi, A\phi \rangle$ if $\psi, \phi \in D(A)$. $q_A$ is called the sesquilinear form associated to $A$, and $Q(q_A)$ is called the form domain of $A$.

Proof - It’s not so clear that $Q(q_A)$ is a vector space; we show this as follows. Suppose $f \in \mathcal{B}(\mathbb{R})$. If $\psi, \phi \in Q(q_A)$,

$$\|f(A)(\psi + \phi)\| \leq \|f(A)\psi\| + \|f(A)\phi\|$$

so using $(a + b)^2 \leq 2(a^2 + b^2)$,

$$\langle (\psi + \phi), f(A)^2(\psi + \phi) \rangle \leq 2\langle \psi, f(A)^2\psi \rangle + 2\langle \phi, f(A)^2\phi \rangle .$$

By the spectral theorem this implies

$$\int f^2 d\mu_{A,\psi + \phi} \leq 2 \int f^2 d\mu_{A,\psi} + 2 \int f^2 d\mu_{A,\phi} .$$

Since $f^2$ can be any positive function in $\mathcal{B}(\mathbb{R})$, this implies that

$$d\mu_{A,\psi + \phi} \leq 2d\mu_{A,\psi} + 2d\mu_{A,\phi} .$$

In particular

$$\int |\lambda| d\mu_{A,\psi + \phi} \leq 2 \int |\lambda| d\mu_{A,\psi} + 2 \int |\lambda| d\mu_{A,\phi} < \infty$$

so $\psi + \phi \in Q(q_A)$.

Given $\psi \in Q(q_A)$, we define $q_A(\psi, \psi)$ as above:

$$q_A(\psi, \psi) = \int_{\sigma(A)} \lambda d\mu_{A,\psi}(\lambda) .$$

Given $\psi, \phi \in Q(q_A)$, we define $q_A(\psi, \phi)$ by the polarization identity; it’s straightforward to verify that this makes $q_A$ really a sesquilinear form. Similarly the polarization identity implies that $q_A$ is the unique form satisfying the conditions of the theorem.
Suppose $\psi \in D(A)$. Then $\psi = (A + i)^{-1}\phi$ for some $\phi \in H$. For all $f \in B(\mathbb{R})$,

\[
\int f \, d\mu_{A,\psi} = \langle \psi, f(A)\psi \rangle = \langle \phi, (A - i)^{-1}f(A)(A + i)^{-1}\phi \rangle = \int \frac{f(\lambda)}{\lambda^2 + 1} \, d\mu_{A,\phi}(\lambda)
\]

so

\[
d\mu_{A,\psi}(\lambda) = (\lambda^2 + 1)^{-1}d\mu_{A,\phi}(\lambda).
\]

In particular,

\[
\int (\lambda^2 + 1) \, d\mu_{A,\psi} = \int 1 \, d\mu_{A,\phi} < \infty,
\]

so

\[
D(A) \subseteq \{ \psi \in H : \int_{\sigma(A)} \lambda^2 \, d\mu_{A,\psi}(\lambda) \} \subseteq Q(q).
\]

Also, for $\psi \in D(A)$

\[
\langle \psi, A\psi \rangle = \int_{\sigma(A)} \lambda \, d\mu_{A,\psi} = q_A(\psi, \psi)
\]

so polarization implies that $q_A(\psi, \phi) = \langle \psi, A\phi \rangle$ if $\psi, \phi \in D(A)$. \hfill \Box

To get a feeling for what’s going on in the theorem, try the following:

**Exercise** - Show that

\[
D(\Delta) = \{ \psi \in H : \int |\Delta \psi|^2 < \infty \}
\]

while

\[
Q(q_\Delta) = \{ \psi \in H : \int |\nabla \psi|^2 < \infty \}.
\]

Given a sesquilinear forms $q$ and $q'$ we write $q \geq q'$ if

\[
q(\psi, \psi) \geq q'(\psi, \psi)
\]

for all $\psi \in Q(q) \cap Q(q')$. We define **bounded below**, **nonnegative**, and **positive** sesquilinear forms by copying the definitions for operators; in particular, $q$ is bounded below if $q$ satisfies $q \geq -cI$ for some $c$. Sesquilinear forms that are bounded below have many nice properties, and for the rest of this section we will discuss forms that are bounded below.
Suppose the sesquilinear form $q$ satisfies $q \geq -cI$. Then there is an inner product on $Q(q)$ given by:

$$\langle \psi, \phi \rangle_1 = q(\psi, \phi) + (c + 1)\langle \psi, \phi \rangle.$$ 

Note that for topological purposes it doesn’t matter which $c$ we pick here:

**Exercise** - Show that if $q \geq -cI$ and $q \geq -c'I$, the two norms on $Q(q)$:

$$\|\psi\|_2^2 = q(\psi, \psi) + (c + 1)\langle \psi, \psi \rangle$$

and

$$\|\psi\|_2' = q(\psi, \psi) + (c' + 1)\langle \psi, \psi \rangle$$

are equivalent, i.e. for some constant $a > 0$,

$$\|\psi\|_1 \leq a\|\psi\|_1', \quad \|\psi\|_1' \leq a\|\psi\|_1.$$ 

Of course $Q(q)$ may not be a Hilbert space under the inner product $\langle \cdot, \cdot \rangle_1$ because it may not be complete in the norm $\|\cdot\|_1$. Given a sesquilinear form $q$ that is bounded below, we define it to be **closed** if $Q(q)$ is complete in the norm $\|\cdot\|_1$. We say that a bounded-below form $q$ is **closable** if there is closed bounded-below form $\overline{q}$ such that $Q(q)$ is dense in $Q(\overline{q})$, the latter given the topology coming from its $\|\cdot\|_1$ norm. In this case we say $\overline{q}$ is the **closure** of $q$:

**Exercise** - Show that the closure of a closable bounded-below sesquilinear form is **unique**.

There is a one-to-one correspondence between closed bounded-below sesquilinear forms and bounded below self-adjoint operators:

**Theorem.** If $A$ is a self-adjoint operator, bounded below, the associated sesquilinear form $q_A$ is closed and bounded below. If $q$ is a sesquilinear form that is closed and bounded below, there is a unique bounded-below self-adjoint operator such that $q = q_A$.

**Proof** - We omit this proof; it can found, for example, in Reed and Simon, *Methods of Modern Mathematical Physics*, vol. I, Theorem VIII.15.

This theorem is rather remarkable because while not all closed operators that are bounded below are self-adjoint, it says that all closed forms that are bounded below correspond to self-adjoint operators. As a consequence, there is a distinguished self-adjoint extension of any operator that is bounded below:
Friedrichs Extension Theorem. Let $A$ be an operator that is bounded below and let $q(\psi, \phi) = \langle \psi, A\phi \rangle$ for $\psi, \phi \in D(A)$. Then $q$ is closable, and the closure $\overline{q}$ equals $q_{\hat{A}}$ for a unique bounded-below self-adjoint operator $\hat{A}$. $\hat{A}$ is an extension of $A$, called the Friedrichs extension of $A$.

Proof - If $A \geq cI$, let $\langle \cdot, \cdot \rangle_1$ be the inner product on $Q(q) = D(A)$ given by:

$$\langle \psi, \phi \rangle_1 = q(\psi, \phi) + (c + 1)\langle \psi, \phi \rangle.$$ 

Let $H_1$ be the Hilbert space completion of $Q(q)$ in the norm $\| \cdot \|_1$. It is easy to see that $q_1$ extends by continuity to a unique form $q_1$ on $H_1$.

Let $\iota: Q(q) \to H$ be the inclusion map. Since $Q(q)$ is dense in $H_1$ and $\| \iota \psi \| \leq \| \psi \|_1$ for $\psi \in Q(q)$, the B.L.T. theorem implies that $\iota$ extends to a unique continuous linear map $\iota_1: H_1 \to H$. We show that $\iota_1$ is one-to-one as follows: if $\iota_1(\psi) = 0$ for some $\psi \in H_1$, then since $Q(q)$ is dense in $H_1$, we can choose $\psi_i \in Q(q)$ such that $\| \psi_i - \psi \|_1 \to 0$,

hence

$$\| \psi_i \| = \| \iota \psi_i \| \to \| \iota_1 \psi \| = 0.$$ 

It follows that

$$\| \psi \|_1^2 = \lim_{j \to \infty} \lim_{i \to \infty} \langle \psi_i, \psi_j \rangle_1$$

$$= \lim_{j \to \infty} \lim_{i \to \infty} \langle \psi_i, A\psi_j \rangle + (c + 1)\langle \psi_i, \psi_j \rangle$$

$$= 0,$$

so $\psi = 0$.

Since $\iota_1$ is one-to-one, it is an isomorphism between $H_1$ and $\text{Ran} \ i_1 \subseteq \mathcal{H}$, allowing us to identify $H_1$ with $\text{Ran} \ i_1$ and $\overline{q}$ with a form on $\text{Ran} \ i_1$, which is clearly the closure $\overline{q}$ since $Q(q)$ is dense in $H_1$ and $H_1$ is complete. Since $\overline{q}$ is closed and bounded below, by the theorem above it is of the form $q_{\hat{A}}$ for a unique bounded-below self-adjoint operator $\hat{A}$. Since for all $\psi, \phi \in D(A)$,

$$\langle \psi, A\phi \rangle = q(\psi, \phi) = \overline{q}(\psi, \phi) = \langle \psi, \hat{A}\phi \rangle,$$

and $D(A)$ is dense in $\mathcal{H}$, we must have $A\phi = \hat{A}\phi$. Thus $\hat{A}$ is an extension of $A$. \[\square\]

An important good feature of the Friedrichs extension is simply its naturality, that is, it takes no extra information to specify which self-adjoint extension of a bounded-below symmetric operator to pick. Thus in the cases where one is trying to get the
correct Hamiltonian out of an operator that is bounded below but not essentially self-adjoint, mathematical physicists take its Friedrichs extension. The following are other nice properties of the Friedrichs extension:

Exercise - Show that if $A$ is a symmetric operator with $A \geq cI$, then $\hat{A} \geq cI$.

Exercise - Show that if $A$ is a symmetric and bounded-below operator then $\hat{A}$ is the only self-adjoint extension of $A$ whose domain is contained in $Q(q)$, where $q$ is as in the theorem.

Exercise - Show that if $A$ is self-adjoint and bounded below then $\hat{A} = A$.

As an example of how the Friedrichs extension can be used on very singular perturbations, we have:

**Corollary.** If $V : \mathbb{R}^3 \to \mathbb{R}$ is continuous except at finitely many points and satisfies $V \geq c$ for some $c \in \mathbb{R}$, then the operator $\Delta + V$ is densely defined and bounded below, hence admits a Friedrichs extension.

Proof - Let $S$ be the set of points on which $V$ is not continuous. Then for $\psi \in C_0^\infty(\mathbb{R}^3 - S)$, $V \psi \in L^2(\mathbb{R}^3)$ and $\Delta \psi \in L^2(\mathbb{R}^3)$. Thus $C_0^\infty(\mathbb{R}^3 - S) \subseteq D(\Delta) \cap D(V)$, so $\Delta + V$ is densely defined. It’s clear also that $\Delta + V$ is bounded below, so the Friedrichs extension exists. □

The above theorem applies to operators of the form $\Delta + cr^{-p}$ with $c \geq 0$.

Here, then, is a summary of results on the Hamiltonian $H_{p,c}$. We haven’t proved all these results, but we’ve proved most of them:

Positive potential - $c > 0$ : $H_{p,c}^0$ is positive hence Friedrichs extension exists.

$0 \leq p < 3/2$ : $H_{p,c}$ essentially self-adjoint on $C_0^\infty(\mathbb{R}^3)$ by Kato-Rellich.

$p = 2, 0 < c \leq 3/4 : H_{p,c}^0$ is not essentially self-adjoint.

$p = 2, c > 3/4 : H_{p,c}^0$ is essentially self-adjoint.

$p > 2 : H_{p,c}^0$ is essentially self-adjoint.

Negative potential - $c < 0$ : $H_{p,c}^0$ is not essentially self-adjoint, but admits self-adjoint extensions.

$0 \leq p < 3/2 : H_{p,c}$ essentially self-adjoint on $C_0^\infty(\mathbb{R}^3)$ by Kato-Rellich.
$0 \leq p < 2 : H^0_{p,c}$ is bounded below hence Friedrichs extension exists.

$p = 2, c \geq -1/4 : H^0_{p,c}$ is bounded below hence Friedrichs extension exists.

$p = 2, c < -1/4 : H^0_{p,c}$ is not bounded below.

$p > 2 : H^0_{p,c}$ is not bounded below.