Unbounded Operators on Hilbert Space with Applications in Quantum Mechanics

Ubegrænsede operatorer på Hilbertrum med anvendelser i kvantemekanik



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Resumé

I denne opgave udvikler vi den basale teori for selvadjungerede operatorer på Hilbertrum og anvender den på adskillige fundamentale eksempler på operatorer som optræder i kvantemekanikken.

Vi antager kendskab til begrænsede operatorer og beviser en spektralsætning for en ubegrænset selvadjungeret operator. Vi betragter stærkt kontinuerte énparameter unitære grupper og beviser Stones sætning som giver en karakterisation af disse grupper. Vi ser da på multiplikationsog differentialoperatorer, beviser selvadjungerethed og beregner spektre for generelle klasser af disse slags operatorer.

Til sidst betragter vi positions- og potentialoperatorer, samt impulsog kinetisk energioperatorer. Vi ser især på sidstnævnte, idet vi udleder tidsudviklingen for en fri partikel.

Abstract

In this thesis we develop the basic theory of self-adjoint operators on Hilbert space and apply it to several fundamental examples of operators appearing in quantum mechanics.

Assuming knowledge of bounded operators, we prove a spectral theorem for an unbounded self-adjoint operator. We consider strongly continuous one-parameter unitary groups and prove Stone's theorem, which gives a characterisation of these groups. We then consider multiplication and differential operators, proving self-adjointness and calculating spectra for general classes of these types of operators.

Finally we consider the position and potential operators, as well as momentum and kinetic energy operators. We pay particular attention to the latter, deriving the time-evolution of a free particle.

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Preface

Until the turn of the twentieth century classical mechanics enjoyed great success, both experimentally and theoretically, the latter in particular in the analytical mechanics due to Lagrange and Hamilton. The end of the nineteenth century saw Planck solve the so-called 'ultraviolet catastrophe' by quantising the exchange of energy in the interaction of matter with electromagnetic radiation, ushering in a new era of physics: quantum mechanics. This development continued in the early 1900s with advances by Einstein, Bohr, Heisenberg and others.

This was paralleled by foundational work in the late 1920s attempting to develop a rigorous mathematical foundation for quantum mechanics. Hilbert spaces and unbounded operators take centre stage, and the systematic development of these operators are first and foremost due to von Neumann (1930) and Stone (1932a).

In this thesis we develop some of this theory. In particular, we prove a spectral theorem for unbounded self-adjoint operators and a theorem due to Stone which characterises certain groups of unitary operators in terms of exponential functions. We then apply these to the Schrödinger equation, after which we consider some of the most important operators in quantum mechanics, prove their self-adjointness and calculate their spectra. We focus on the kinetic energy operator, deriving the time-evolution of a free particle. An attempt has been made to be as general as possible without detracting from the main goal, which is the application of these theorems in quantum mechanics.

We assume knowledge of measure theory and functional analysis as seen in the course Advanced Analysis at Aarhus University, roughly at the level of Rudin (1991). In particular, it is assumed that the reader knows the basic theory of bounded operators on Hilbert space and fundamental theorems such as the Hahn–Banach theorem. Prerequisites not explicitly covered in this course or bachelor courses at Aarhus University, such as the *d*-dimensional Fourier transform and the Baire category theorem, are collected in an appendix. We either prove all claims or explicitly refer to the literature. A passing familiarity with quantum mechanics would also be beneficial to appreciate the material,

Preface

but should not be strictly necessary.

For the general theory of unbounded operators the main references have been Rudin (1991), Schmüdgen (2012) and the notes Skibsted (2003). For applications in quantum mechanics as well as historical notes I have greatly benefited from Reed and Simon (1975) and Hall (2013). Sakurai and Napolitano (2011) has been my go-to for the physicist's perspective. The texts Rudin (1986) and in particular Folland (2007) have also been very useful as general references.

It is a pleasure to give a thanks first of all to Erik Skibsted for an inspiring course in Advanced Analysis, and for his effective guidance while writing this thesis. Also a thank you to Steen Thorbjørnsen for teaching me analysis, and to Søren Fournais for introducing me to mathematical physics.

The layout and typography of the present thesis were implemented by myself, but much has been borrowed from Steen Thorbjørnsen's wonderful (both in content and in style) *Grundlæggende mål- og integralteori*. I also owe a great debt to Lars Madsen, whose *Introduktion til LATEX* has been a tremendous help and source of inspiration over the years.

Definitions and theorems are placed in coloured boxes to set them apart from the surrounding text. Note that these boxes may span multiple pages. Following common practice, we use the tombstone ' \Box ' to denote the end of a proof. Examples and remarks are terminated with a ' \Box '.

Aarhus, January 2021

Danny Nygård Hansen

— 1 —

Unbounded operators

1.1 • Definitions and examples

DEFINITION 1.1

If *X* and *Y* are vector spaces, an *operator* from *X* to *Y* is a linear map *T* from a subspace of *X* into *Y*. This subspace is called the *domain* of *T* and is denoted $\mathcal{D}(T)$. If *T* is an operator from *X* to *Y*, we also write $T: X \to Y$.

The spaces *X* and *Y* will always be normed, and to emphasise that an operator *T* is not necessarily bounded, we will sometimes call such an operator *unbounded*, although a better name might have been *not necessarily bounded*.

Our motivation for introducing unbounded operators is that many of the operators appearing in quantum mechanics (and physics more broadly) are unbounded. An important class of operators is the class of differential operators, illustrated by the following simple example:

EXAMPLE 1.2: An unbounded differential operator.

Equip C[0,1] with the supremum norm, and let $T: C[0,1] \to C[0,1]$ be the differential operator Tf = f' with domain $\mathcal{D}(T) = C^1[0,1]$. This is clearly well-defined and linear. But consider the sequence of functions (f_n) given by $f_n(x) = \mathrm{e}^{\mathrm{i}\, nx}$. Then $\|f_n\|_{\infty} = 1$, but $\|Tf_n\|_{\infty} = n$, so T is not bounded on its domain.

The operator in the example above is not defined everywhere on C[0,1]. However, there do exist operators that are everywhere-defined but not bounded, as the following example shows:

EXAMPLE 1.3: An everywhere-defined not bounded operator.

Let X be an infinite-dimensional normed space, and let $Y \neq 0$ be a normed space. Choose a Hamel basis \mathcal{B} for X, and let $\mathcal{B}' = \{b_n \mid n \in \mathbb{N}\}$ be a countable subset of \mathcal{B} . Pick some $y \in Y$ with $y \neq 0$. Then define

$$T(b) = \begin{cases} n||b_n||y, & b = b_n, \\ 0, & b \in \mathcal{B} \setminus \mathcal{B}', \end{cases}$$

and extend T to an operator on X. Then $T: X \to Y$ is clearly everywhere-defined and not bounded.

If $S,T:X\to Y$ are operators with $\mathcal{D}(S)\subseteq\mathcal{D}(T)$ such that Sx=Tx for all $x\in\mathcal{D}(S)$, then we write $S\subseteq T$ and say that T is an *extension* of S. This is clearly equivalent to the condition $\mathcal{G}(S)\subseteq\mathcal{G}(T)$, where $\mathcal{G}(T)$ denotes the graph of an operator T.

Since operators are not necessarily defined on the whole space, one must take care when e.g. adding or composing operators. We refer the reader to Rudin (1991, Definitions 13.1) or Schmüdgen (2012, Section 1.1.1) for details. Furthermore, note that for an operator $T: \mathcal{H} \to \mathcal{K}$ between Hilbert spaces that is bounded on $\mathcal{D}(T)$, we can always extend it to \mathcal{H} : First extend it to $\overline{\mathcal{D}(T)}$ using the BLT Theorem B.2, and then let Tx = 0 for $x \in \overline{\mathcal{D}(T)}^{\perp}$.

An important class of operators are the *closed operators*:

DEFINITION 1.4: Closed and closable operators

An operator $T: X \to Y$ between normed spaces is *closed* if its graph $\mathcal{G}(T)$ is a closed subset of the direct sum $X \oplus Y$. Furthermore, T is *closable* if the closure of its graph is the graph of some function. The function corresponding to $\overline{\mathcal{G}(T)}$ is called the *closure* of T and is denoted \overline{T} .

If $\overline{\mathcal{G}(T)}$ is the graph of a function, that function is clearly an operator $\overline{T} \colon X \to Y$ since $\overline{\mathcal{G}(T)}$ is a subspace of $X \oplus Y$. If X and Y are Banach spaces, it follows from the Closed Graph Theorem B.5 that $T \in \mathcal{B}(X,Y)$ if and only if $\mathcal{D}(T) = X$ and T is closed.

DEFINITION 1.5: Spectrum and resolvent of an operator

An operator $T \colon X \to Y$ between Banach spaces is *boundedly invertible* if there is a bounded operator $S \colon Y \to X$ such that TS = 1 and $ST \subseteq 1$. The *resolvent set* $\rho(T)$ of T is the set of complex numbers λ such that $\lambda - T$ is boundedly invertible. The *spectrum* of T is the set $\sigma(T) = \mathbb{C} \setminus \rho(T)$.

Note that the Bounded Inverse Theorem B.4 ensures that this definition of the spectrum agrees with that of a bounded operator if $T \in \mathcal{B}(X, Y)$.

In the sequel, we will restrict to operators between Hilbert spaces, and the symbols \mathcal{H} and \mathcal{K} will denote arbitrary Hilbert spaces.

1.2 • Symmetric and self-adjoint operators

Let $T: \mathcal{H} \to \mathcal{K}$ be an operator. We wish to associate an adjoint operator $T^*: \mathcal{K} \to \mathcal{H}$ to T. First let $\mathcal{D}(T^*)$ be the set of all $y \in \mathcal{K}$ such that the linear

¹ More generally from a Hilbert space into a Banach space.

functional $x \mapsto \langle Tx, y \rangle$ is continuous on $\mathcal{D}(T)$. If $y \in \mathcal{D}(T^*)$, then the Hahn–Banach theorem extends the above functional to a continuous linear functional $\varphi \colon \mathcal{H} \to \mathbb{C}$. The Riesz–Fréchet theorem then yields an element $y_{\varphi} \in \mathcal{H}$ such that $\varphi(x) = \langle x, y_{\varphi} \rangle$. In particular, for $x \in \mathcal{D}(T)$ we have $\langle Tx, y \rangle = \langle x, y_{\varphi} \rangle$. If T is densely defined, then clearly y_{φ} is unique. Otherwise it is not, since for every $z \in \mathcal{D}(T)^{\perp} \setminus \{0\}$ we have $\langle x, y_{\varphi} \rangle = \langle x, y_{\varphi} + z \rangle$.

If T is densely defined we let $T^*y = y_{\varphi}$. Then it is clear that $T^* \colon \mathcal{K} \to \mathcal{H}$ is an operator with domain $\mathcal{D}(T^*)$. The operator T^* is called the *adjoint* of T. Do note that T^* is not in general densely defined, but is so if and only if T is closable, cf. Proposition 1.10(iii) below.

Like bounded operators, unbounded operators can also be *self-adjoint*. However, we must take care to distinguish this from the weaker property of *symmetry*:

DEFINITION 1.6: Symmetric and self-adjoint operators

An operator T on \mathcal{H} is *symmetric* if

$$\langle Tx, y \rangle = \langle x, Ty \rangle$$

for all $x, y \in \mathcal{D}(T)$. If T is densely defined and $T = T^*$, then we say that T is self-adjoint. If T is densely defined and closable, we say that it is essentially self-adjoint if \overline{T} is self-adjoint.

The densely defined symmetric operators are then precisely those operators that satisfy $T \subseteq T^*$.

A symmetric operator T on \mathcal{H} is called *maximally symmetric* if it has no proper symmetric extensions, i.e. if $T \subseteq S$ with S symmetric implies T = S. If T is self-adjoint, then it is maximally symmetric: For if $T \subseteq S$ with S symmetric, then $S^* \subseteq T^*$. But $T \subseteq S \subseteq S^* \subseteq T^* = T$, so T = S.

EXAMPLE 1.7: The position and momentum operators.

Define operators X and P on $L^2(\mathbb{R})$ by Xf(x) = xf(x) and $Pf = -\mathrm{i}\,f'$ on the natural domains

$$\begin{split} \mathcal{D}(X) &= \{ f \in L^2(\mathbb{R}) \mid \mathrm{id}_{\mathbb{R}} f \in L^2(\mathbb{R}) \}, \\ \mathcal{D}(P) &= \{ f \in L^2(\mathbb{R}) \mid f \in C_0^1(\mathbb{R}), f' \in L^2(\mathbb{R}) \}. \end{split}$$

Then it is clear that X is symmetric. That P is also symmetric follows easily from integration by parts since the boundary term vanishes. Since $C_c^{\infty}(\mathbb{R})$ is contained in each domain, both operators are densely defined by Proposition B.1. We will study these and other operators from quantum mechanics in more detail in Section 5.

² Alternatively, the BLT Theorem B.2 extends the functional to a bounded linear functional on the closed subspace $\overline{\mathcal{D}(T)}$, and we may then extend it to \mathcal{H} using the projection theorem.

LEMMA 1.8

If $T: \mathcal{H} \to \mathcal{K}$ *is a densely defined operator, then*

$$\mathcal{R}(T)^{\perp} = \mathcal{N}(T^*).$$

PROOF. First let $x \in \mathcal{R}(T)^{\perp}$. Then $\langle Ty, x \rangle = 0$ for all $y \in \mathcal{D}(T)$, and hence $x \in \mathcal{D}(T^*)$. It follows that $\langle y, T^*x \rangle = 0$, and since T is densely defined, this implies that $T^*x = 0$.

Conversely, if $x \in \mathcal{N}(T^*)$, then $0 = \langle y, T^*x \rangle = \langle Ty, x \rangle$ for all $y \in \mathcal{D}(T)$, showing that $x \in \mathcal{R}(T)^{\perp}$.

We now prove a series of important results about adjoints and symmetric operators. They are proved using a useful technique known as the *graph method*, which is due to von Neumann (1932). We present this method in the following lemma. It is possible to prove several other interesting results using the graph method, see in particular Rudin (1991, Theorem 13.11) and Schmüdgen (2012, Theorem 1.8).

LEMMA 1.9

Let $T: \mathcal{H} \to \mathcal{K}$ be densely defined, and define an operator $V: \mathcal{H} \oplus \mathcal{K} \to \mathcal{K} \oplus \mathcal{H}$ by V(x,y) = (-y,x). Then

$$\mathcal{G}(T^*) = (V\mathcal{G}(T))^{\perp} = V\mathcal{G}(T)^{\perp}.$$

PROOF. We have $(y,z) \in \mathcal{G}(T^*)$ iff $\langle Tx,y \rangle = \langle x,z \rangle$ for all $x \in \mathcal{D}(T)$ by definition of T^* . But this is true iff $\langle (-Tx,x),(y,z) \rangle = 0$ for all $x \in \mathcal{D}(T)$, which is equivalent to $(y,z) \in (V\mathcal{G}(T))^{\perp}$. The second equality follows since V is unitary. \square

PROPOSITION 1.10

Let $T: \mathcal{H} \to \mathcal{K}$ be densely defined. Then:

- (i) T^* is closed. In particular, if T is symmetric then it is closable, and if it is self-adjoint, then it is closed.
- (ii) If T is closable, then $(\overline{T})^* = T^*$.
- (iii) T is closable iff T^* is densely defined.
- (iv) If T is closable, then $\overline{T} = T^{**}$.

PROOF. We begin with (i). The first claim is obvious from Lemma 1.9, and the second follows since T^* is a closed extension of T if T is symmetric.

Assume that *T* is closable and note that

$$\mathcal{G}((\overline{T})^*) = V\mathcal{G}(\overline{T})^{\perp} = V\overline{\mathcal{G}(T})^{\perp} = V\mathcal{G}(T)^{\perp} = \mathcal{G}(T^*).$$

This proves (ii). For $y \in \mathcal{D}(T^*)^{\perp}$ we have

$$(-v,0) \in \mathcal{G}(T^*)^{\perp} = (V\mathcal{G}(T))^{\perp \perp} = V\overline{\mathcal{G}(T)} = V\overline{\mathcal{G}(T)},$$

which implies that $(0, y) \in \overline{\mathcal{G}(T)}$. But if T is closable, then $\overline{\mathcal{G}(T)}$ is the graph of an operator, so y = 0. Conversely, assume that T^* is densely defined. Then T^{**} is closed by (i), and we clearly have $T \subseteq T^{**}$. Thus T has a closed extension and is thus closable. This proves (iii).

Again assume that T is closable. By (iii), T^* is densely defined, so T^{**} exists. Notice that the operator $-V^{-1}$ corresponds to T^* in the same way that V corresponds to T. Then

$$\mathcal{G}(T^{**}) = -V^{-1}\mathcal{G}(T^{*})^{\perp} = V^{-1}V\mathcal{G}(T)^{\perp \perp} = \overline{\mathcal{G}(T)} = \mathcal{G}(\overline{T}),$$

proving (iv).

COROLLARY 1.11: The Hellinger-Toeplitz Theorem

If T is a symmetric operator on \mathcal{H} with $\mathcal{D}(T) = \mathcal{H}$, then T is bounded.

PROOF. It follows from Proposition 1.10(i) that T is closable, but $T = \overline{T}$ on $\mathcal{D}(T) = \mathcal{D}(\overline{T})$, so T is already closed. Since it is defined everywhere, the Closed Graph Theorem B.5 implies that it is bounded.

In Example 1.7 we saw that the position and momentum operators are symmetric on natural domains. But they are clearly not bounded, so the Hellinger–Toeplitz theorem implies that we cannot hope to find symmetric (and hence neither self-adjoint) extensions to \mathcal{H} . Instead we will have to settle for them being self-adjoint on dense domains.

Operators on the form $T \pm i$ will be important in the definition of the Cayley transform. We end this section by proving some properties of these operators that we will need.

PROPOSITION 1.12

Let T be a symmetric operator on \mathcal{H} . Then the following statements hold:

- (i) $||Tx + ix||^2 = ||x||^2 + ||Tx||^2$ for all $x \in \mathcal{D}(T)$.
- (ii) T is closed if and only if R(T + i) is closed.
- (iii) T + i is injective.

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The above statements also hold with i replaced by -i.

PROOF. Since *T* is symmetric, the calculation

$$||Tx + ix||^2 = ||Tx||^2 + ||x||^2 + \langle Tx, ix \rangle + \langle ix, Tx \rangle$$

proves (i). Now notice that the map $\mathcal{G}(T) \to \mathcal{R}(T+i)$ given by $(x, Tx) \mapsto (T+i)x$ is isometric by (i). It is also clearly surjective, so this proves (ii). The claim (iii) also follows immediately since if (T+i)x = (T+i)y, then (x, Tx) = (y, Ty).

The final claim follows by replacing T with -T.

1.3 • Conditions for self-adjointness

We now prove a criterion for self-adjointness that will be useful in the development of the Cayley transform.

PROPOSITION 1.13

Let T be a symmetric operator on \mathcal{H} . If there exists a $\lambda \in \mathbb{C}$ such that $\mathcal{R}(T - \lambda I) = \mathcal{H}$ and $\mathcal{R}(T - \overline{\lambda}I)$ is dense in \mathcal{H} , then T is self-adjoint.

PROOF. We first show that T is densely defined. By the projection theorem it is enough to show that $\mathcal{D}(T)^{\perp} = 0$. So let $y \in \mathcal{D}(T)^{\perp}$ and write $y = (T - \lambda I)x$ for some $x \in \mathcal{D}(T)$. For all $z \in \mathcal{D}(T)$ we have

$$0 = \langle y, z \rangle = \langle (T - \lambda I)x, z \rangle = \langle x, (T - \overline{\lambda}I)z \rangle,$$

and since $\mathcal{R}(T - \overline{\lambda}I)$ is dense in \mathcal{H} , this implies that y = 0.

Since T is densely defined, its adjoint T^* exists. Since T is symmetric, $T \subseteq T^*$, so it is enough to show that $\mathcal{D}(T^*) \subseteq \mathcal{D}(T)$. Let $y \in \mathcal{D}(T^*)$ and let $z \in \mathcal{D}(T)$ be a vector such that $(T^* - \lambda I)y = (T - \lambda I)z$. For all $x \in \mathcal{D}(T)$ we have

$$\langle (T - \overline{\lambda}I)x, y \rangle = \langle x, (T^* - \lambda I)y \rangle = \langle x, (T - \lambda I)z \rangle = \langle (T - \overline{\lambda}I)x, z \rangle.$$

Again since $\mathcal{R}(T - \overline{\lambda}I)$ is dense in \mathcal{H} , we have y = z, so $y \in \mathcal{D}(T)$.

DEFINITION 1.14

Let T be an operator on \mathcal{H} . The *deficiency spaces* of T are the subspaces $\mathcal{L}_{\pm} = \mathcal{R}(T \pm i)^{\perp}$, and the *deficiency indices* are the cardinal numbers³ $n_{\pm} = \dim \mathcal{L}_{\pm}$.

Note that some authors use opposite signs in the definition of \mathcal{L}_{\pm} and n_{\pm} .

³ Note that \mathcal{L}_{\pm} are closed in \mathcal{H} , so they are themselves Hilbert spaces. By definition, the dimension of a Hilbert space is the cardinality of any of its orthonormal bases.

REMARK 1.15. If T is a closable and densely defined operator, then the deficiency spaces of T and \overline{T} are the same. For by Proposition 1.10(ii) and Lemma 1.8 we have

$$\mathcal{R}(T \pm i)^{\perp} = \mathcal{N}(T^* \mp i) = \mathcal{N}((\overline{T})^* \mp i) = \mathcal{R}(\overline{T} \pm i)^{\perp}.$$

PROPOSITION 1.16

Let T be a densely defined symmetric operator on \mathcal{H} . Then T is essentially self-adjoint if and only if its deficiency indices satisfy $n_{\pm}=0$.

PROOF. First assume that T is also closed. If T is self-adjoint, then $T^*\pm i = T\pm i$ is injective by Proposition 1.12(iii), and so $n_{\pm} = 0$ by Lemma 1.8. Conversely, if $n_{\pm} = 0$, then $\mathcal{R}(T \pm i)^{\perp} = 0$ by definition. And since $\mathcal{R}(T \pm i)$ is closed by Proposition 1.12(ii), $\mathcal{R}(T \pm i) = \mathcal{H}$, so T is self-adjoint by Proposition 1.13.

If T is not assumed to be closed, the statement follows by applying the above argument to \overline{T} (cf. Remark 1.15).

_ 2 _

The Cayley transform

The Möbius transformation

$$t \mapsto \frac{t - i}{t + i} \tag{2.1}$$

is a bijection from the real line to the set $S^1 \setminus \{1\}$, i.e. the unit circle in \mathbb{C} without the point 1. The Cayley transform is an analogue of this transformation for operators, specifically between symmetric operators and a certain class of isometries.⁴

2.1 • Properties of the Cayley transform

We begin with the definition:

⁴ We remark that it is also possible to define a Cayley transform based on the Möbius transformation $t \mapsto (t-\lambda)(t-\overline{\lambda})$ for any complex number λ with Im $\lambda > 0$. See Schmüdgen (2012, Section 13.1) for such a development for densely defined symmetric operators.

DEFINITION 2.1: The Cayley transform of a symmetric operator

Let *T* be a symmetric operator on \mathcal{H} . The *Cayley transform* of *T* is the isometry *U* with domain $\mathcal{D}(U) = \mathcal{R}(T + i)$ given by

$$U = (T - i)(T + i)^{-1}$$
.

The range of *U* is $\mathcal{R}(U) = \mathcal{R}(T - i)$.

From Proposition 1.12(iii) we know that T + i is injective, so $(T + i)^{-1}$ exists. Since $(T + i)^{-1}$ maps $\mathcal{D}(U) = \mathcal{R}(T + i)$ onto $\mathcal{D}(T + i) = \mathcal{D}(T - i)$, the domain of U makes sense, and its range is in fact $\mathcal{R}(T - i)$.

Any element of $\mathcal{R}(T+i)$ is on the form Tx+ix for some $x \in \mathcal{D}(T)$, so we can also write

$$U(Tx + ix) = Tx - ix$$
.

It then follows from Proposition 1.12(i) that *U* is an isometry as claimed.

THEOREM 2.2: The inverse Cayley transform

Let T be a symmetric operator on \mathcal{H} with Cayley transform U. Then the operator I-U is injective with range $\mathcal{R}(I-U)=\mathcal{D}(T)$. Furthermore,

$$T = i(I + U)(I - U)^{-1}.$$
 (2.2)

In particular, different symmetric operators have different Cayley transforms.

Conversely, every isometry U on \mathcal{H} such that I-U is injective is the Cayley transform of some symmetric operator on \mathcal{H} .

Before proving this theorem, we remark that the inverse of the Möbius transformation (2.1) is

$$z \mapsto i \frac{1+z}{1-z}$$
,

which explains the formula (2.2). Also note that the theorem in particular says that the Cayley transform $T \mapsto (T-i)(T+i)^{-1}$ gives a bijective correspondence between the set of symmetric operators and the set of those isometries U such that I-U is injective, and that the *inverse Cayley transform* is given by $U \mapsto i(I+U)(I-U)^{-1}$.

PROOF OF THEOREM 2.2. Let $z \in \mathcal{D}(U) = \mathcal{R}(T+\mathrm{i})$. Since $T+\mathrm{i}$ is injective, it defines a bijection between $\mathcal{D}(T)$ and $\mathcal{D}(U)$, so there is a unique $x \in \mathcal{D}(T)$ such that $z = Tx + \mathrm{i} x$. Applying U we get $Uz = Tx - \mathrm{i} x$, and combining these two formulas yields

$$(I-U)z = 2ix$$
 and $(I+U)z = 2Tx$.

A $z' \neq z$ would yield an $x' \neq x$ such that (I - U)z' = 2ix', showing that I - U is injective. The first identity also shows that $\mathcal{R}(I - U) = \mathcal{D}(T)$, so $(I - U)^{-1}$ maps $\mathcal{D}(T)$ onto $\mathcal{D}(U)$. Using the second identity we get

$$2Tx = (I + U)z = (I + U)(I - U)^{-1}(2ix),$$

proving (2.2).

We now prove the converse statement. Let U be an isometry such that I - U is injective and define an operator T with domain $\mathcal{D}(T) = \mathcal{R}(I - U)$ by

$$T = i(I + U)(I - U)^{-1}$$
.

We first show that T is symmetric. For $x, y \in \mathcal{D}(T)$ write x = (I - U)z and y = (I - U)w for appropriate $z, w \in \mathcal{D}(U)$. Then using the fact that U preserves inner products, we get

$$\langle Tx, y \rangle = i \langle (I+U)z, (I-U)w \rangle = i \langle Uz, w \rangle - i \langle z, Uw \rangle$$

= $-i \langle (I-U)z, (I+U)w \rangle = \langle x, Ty \rangle$,

as desired. Since T(I-U)z = i(I+U)z, it is a quick calculation to show that

$$(T+i)(I-U)z = 2iz$$
 and $(T-i)(I-U)z = 2iUz$.

Because T is symmetric, T + i is injective by Proposition 1.12(iii), so it follows that

$$2i Uz = (T - i)(T + i)^{-1}(2iz),$$

showing that $U \subseteq (T - i)(T + i)^{-1}$. On the other hand, we have

$$\mathcal{D}((T-\mathrm{i})(T+\mathrm{i})^{-1}) \subseteq \mathcal{D}((T+\mathrm{i})^{-1}) = \mathcal{R}(T+\mathrm{i}) = \mathcal{D}(U),$$

so in total *U* is the Cayley transform of *T*.

REMARK 2.3. If T is a densely defined symmetric operator with Cayley transform U, then $\mathcal{R}(I-U)=\mathcal{D}(T)$ is dense in \mathcal{H} by Theorem 2.2. Conversely, if U is an isometry on \mathcal{H} such that $\mathcal{R}(I-U)$ is dense in \mathcal{H} , then I-U is injective. To see this, let $x \in \mathcal{N}(I-U)$ and $y \in \mathcal{D}(I-U)$. Then x = Ux, so

$$\langle x, (I-U)y \rangle = \langle x, y \rangle - \langle x, Uy \rangle = \langle x, y \rangle - \langle Ux, Uy \rangle = \langle x, y \rangle - \langle x, y \rangle = 0,$$

and since $\mathcal{R}(I-U)$ is dense, x=0. It follows by Theorem 2.2 that U has an inverse Cayley transform T which is densely defined since $\mathcal{D}(T) = \mathcal{R}(I-U)$.

The Cayley transform thus gives a one-to-one correspondence between densely defined symmetric operators and those isometries U such that $\mathcal{R}(I-U)$ is dense in \mathcal{H} .

We now present the form of the Cayley transform that we will use in the proof of the spectral theorem:

THEOREM 2.4: The Cayley transform for self-adjoint operators

Let T be a symmetric operator on H with Cayley transform U. Then T is self-adjoint if and only if U is unitary.

PROOF. First assume that T is self-adjoint. Then $T \pm i$ is injective by Proposition 1.12(iii), so Lemma 1.8 implies that

$$\mathcal{R}(T \pm i)^{\perp} = \mathcal{N}(T \mp i) = 0,$$

and since T is closed by Proposition 1.10(i), it follows from Proposition 1.12(ii) that $\mathcal{R}(T\pm i)=\mathcal{H}$. Then Proposition 1.13 yields that T is self-adjoint if and only if $\mathcal{R}(T\pm i)=\mathcal{H}$. But $\mathcal{D}(U)=\mathcal{R}(T+i)$ and $\mathcal{R}(U)=\mathcal{R}(T-i)$, so T is self-adjoint if and only if U is unitary.

2.2 • Extensions of symmetric operators

The Cayley transform plays a crucial role in our proof of the spectral theorem, but we now take a brief detour to discuss another application: The Cayley transform can also be used to characterise symmetric and self-adjoint extensions of a densely defined symmetric operator T in terms of respectively isometric and unitary extensions of the Cayley transform of U. We will not need this result in the sequel, but it is an important development in the theory of self-adjoint extensions of symmetric operators. We have already seen in Example 1.7 that it might be easy to write down domains on which operators are symmetric, but self-adjointness is more difficult as we will see in Section 5.

The approach to self-adjoint extension theory using the Cayley transform is due to von Neumann, and we refer to Schmüdgen (2012, Part VI) for further discussion.

COROLLARY 2.5: Characterisation of symmetric extensions

Let T be a densely defined symmetric operator on $\mathcal H$ with Cayley transform U.

- (i) The Cayley transform gives a one-to-one correspondence between symmetric extensions of T and isometric extensions of U.
- (ii) The self-adjoint extensions of T correspond exactly to the unitary extensions of U.

PROOF. First let \tilde{T} be a symmetric operator with Cayley transform \tilde{U} . To prove (i) it is enough by Remark 2.3 to show that $T \subseteq \tilde{T}$ if and only if $U \subseteq \tilde{U}$. If $T \subseteq \tilde{T}$, then $\mathcal{D}(U) = \mathcal{R}(T+i) \subseteq \mathcal{R}(\tilde{T}+i) = \mathcal{D}(\tilde{U})$. The converse follows similarly by using Theorem 2.2.

Claim (ii) now follows directly from Theorem 2.4.

THEOREM 2.6: Characterisation of self-adjoint extensions

Let T be a closed and densely defined symmetric operator on \mathcal{H} with deficiency spaces \mathcal{L}_{\pm} . The symmetric (resp. self-adjoint) extensions of T are then in a one-to-one correspondence with the set of isometries (resp. unitary operators) $\mathcal{L}_{+} \to \mathcal{L}_{-}$.

In particular, T has a self-adjoint extension if and only if $n_+ = n_-$, and this extension is unique if and only if $n_\pm = 0$, in which case T is already self-adjoint.

PROOF. First recall that $\mathcal{R}(T \pm i)$ is closed by Proposition 1.12(ii), and notice that the projection theorem yields the decompositions

$$\mathcal{H} = \mathcal{R}(T \pm i) \oplus \mathcal{R}(T \pm i)^{\perp} = \mathcal{R}(T \pm i) \oplus \mathcal{L}_{\pm}.$$

If U is the Cayley transform of T, then $\mathcal{D}(U) = \mathcal{R}(T+\mathrm{i})$ and $\mathcal{R}(U) = \mathcal{R}(T-\mathrm{i})$, and so an isometric extension of U corresponds to an isometry $\mathcal{L}_+ \to \mathcal{L}_-$. Furthermore, if such an isometry is to be unitary, then we clearly need $n_+ = n_-$. Conversely, there are many such unitary operators unless $n_\pm = 0$, in which case there is only the zero operator. The result then follows from Corollary 2.5. \square

— 3 **—**

The spectral theorem

There are many versions of the spectral theorem for self-adjoint operators on a Hilbert space, both for bounded and unbounded operators. The version we present here uses spectral measures⁵, and the presentation is mainly based on Rudin (1991, Chapter X) and Skibsted (2003, Section 4). Our proof is via the Cayley transform which allows us to reduce the problem to the case of a bounded normal operator. For an approach that does not assume the spectral theorem for bounded normal operators, but instead uses operator Riemann–Stieltjes integrals, see Schmüdgen (2012, Chapter 4). Also see Appendix C.1 for further comments about spectral measures.

For reference, Reed and Simon (1980) and Hall (2013) provide several other versions of the spectral theorem for separable Hilbert spaces. Notice that we do not assume that our Hilbert spaces are separable in our presentation.

⁵ These objects have various other names, such as 'resolutions of the identity' and 'projection-valued measures'.

It is also possible to prove spectral theorems for *normal* unbounded operators: An closed and densely defined operator T on a Hilbert space \mathcal{H} is normal if $TT^* = T^*T$, analogous to the bounded case. However, we shall have no need for such a theorem. See Rudin (1991, Theorem 13.33) or Conway (1990, Theorem 4.11) for the case of a single normal operator, and see Schmüdgen (2012, Theorem 5.21) for a spectral theorem for finitely many so-called *strongly commuting* normal operators.

3.1 • Spectral measures

Throughout this section (Ω, \mathcal{E}) will denote a measurable space, and \mathcal{H} a nonzero Hilbert space. We assume that the reader is familiar with spectral measures and refer to Appendix C.1 for definitions and notation.

Our first task will be to define the integral of a measurable function with respect to a spectral measure. Below we let E denote a fixed spectral measure. If $f \in L^{\infty}(E)$, then Theorem C.3 shows that the resulting integral $\int f \, \mathrm{d}E$ is a bounded operator on \mathcal{H} . If f is not essentially bounded, then the corresponding integral $\int f \, \mathrm{d}E$ is also not bounded, as we will see in Proposition 3.8. Its domain will turn out to be the set \mathcal{D}_f in the following lemma:

LEMMA 3.1

For all $f \in \mathcal{M}(\mathcal{E})$ the set

$$\mathcal{D}_f = \{ x \in \mathcal{H} \mid f \in L^2(E_x) \} = \left\{ x \in \mathcal{H} \mid \int |f|^2 \, \mathrm{d}E_x < \infty \right\}$$
 (3.1)

is a dense subspace of H.

Before proving Lemma 3.1, note that for $a, b \in \mathbb{R}$ we have the inequality

$$(a+b)^2 \le 2a^2 + 2b^2. \tag{3.2}$$

This follows e.g. from the calculation

$$2a^2 + 2b^2 - (a+b)^2 = a^2 + b^2 - 2ab = (a-b)^2 \ge 0.$$

PROOF OF LEMMA 3.1. We first show that \mathcal{D}_f is a subspace of \mathcal{H} , so let $x, y \in \mathcal{D}_f$, $A \in \mathcal{E}$ and put z = x + y. Since E(A) is an orthogonal projection, we have $E_x(A) = ||E(A)x||^2$ so it follows by (3.2) that

$$E_z(A) = ||E(A)z||^2 \le 2||E(A)x||^2 + 2||E(A)y||^2 = 2E_z(A) + 2E_z(A).$$

Thus $z \in \mathcal{D}_f$. Closure under scalar multiplication follows similarly.

Next we show that \mathcal{D}_f is dense in \mathcal{H} . First let A_n be the (measurable) set on which $|f| \le n$ for $n \in \mathbb{N}$. If $x \in \mathcal{R}(E(A_n))$, then $x = E(A_n)x$ since $E(A_n)$ is a projection. Hence

$$E(A)x = E(A)E(A_n)x = E(A \cap A_n)x$$

for all $A \in \mathcal{E}$, so $E_x(A) = E_x(A \cap A_n)$. But then we clearly have

$$\int_{\Omega} |f|^2 dE_x = \int_{A_n} |f|^2 dE_x \le n^2 ||x||^2 < \infty,$$

yielding $x \in \mathcal{D}_f$. This shows that $\mathcal{R}(E(A_n)) \subseteq \mathcal{D}_f$. Now let $y \in \mathcal{H}$. We wish to find a sequence $(y_n)_{n \in \mathbb{N}}$ with $y_n \in \mathcal{R}(E(A_n)) \subseteq \mathcal{D}_f$ that converges to y. But $\Omega = \bigcup_{n=1}^{\infty} A_n$, so

$$y = E(\Omega)y = E\left(\bigcup_{n=1}^{\infty} A_n\right)y = \lim_{n \to \infty} E(A_n)y.$$

Thus y lies in the closure of \mathcal{D}_f , and so \mathcal{D}_f is dense in \mathcal{H} .

THEOREM 3.2: Integration with respect to a spectral measure

For every $f \in \mathcal{M}(\mathcal{E})$ there is a densely defined operator $\Psi(f)$ in \mathcal{H} with domain $\mathcal{D}(\Psi(f)) = \mathcal{D}_f$ such that

$$\langle \Psi(f)x, x \rangle = \int_{\Omega} f \, dE_x, \quad (x \in \mathcal{D}_f),$$
 (3.3)

and

$$\|\Psi(f)x\|^2 = \int_{\Omega} |f|^2 dE_x, \quad (x \in \mathcal{D}_f).$$
 (3.4)

The operator $\Psi(f)$ is uniquely determined by its domain and (3.3). We call $\Psi(f)$ the spectral integral of f with respect to E.

In the sequel $\Psi(f)$ will always denote the spectral integral of a measurable function f. The spectral measure E will usually be obvious from context. We will also use the notation $\int f \, dE$ for this operator.

PROOF OF THEOREM 3.2. This proof is similar to that of Skibsted (2019, Theorem 3.56), so we skip some details throughout. Notice that, since E_x is a finite measure for all $x \in \mathcal{H}$, we have $L^2(E_x) \subseteq L^1(E_x)$, so the integral in (3.3) makes sense.

For $s \in \mathcal{SM}(\mathcal{E})$ with standard representation $s = \sum_{j=1}^{n} a_j \mathbf{1}_{A_j}$, notice that

$$\left\langle \sum_{j=1}^{n} a_j E(A_j) x, x \right\rangle = \sum_{j=1}^{n} a_j E_x(A_j) = \int s \, \mathrm{d}E_x.$$

By polarisation, the map Ψ defined by

$$\Psi(s) = \sum_{j=1}^{n} a_j E(A_j).$$

is then well-defined, and Ψ satisfies (3.3) with s in place of f. The map Ψ is also linear on $\mathcal{SM}(\mathcal{E})$, again by polarisation.

We now extend Ψ to all measurable functions. The identity (3.4) for all $x \in \mathcal{H}$ follows with s in place of f by the multiplicativity of E (cf. Proposition C.2). In particular, for a sequence $(s_n) \subseteq \mathcal{SM}(\mathcal{E})$ we have

$$\|\Psi(s_n)x - \Psi(s_m)x\|^2 = \int |s_n - s_m|^2 dE_x$$
 (3.5)

for $x \in \mathcal{H}$.

Let $f \in \mathcal{M}(\mathcal{E})$, and let (s_n) be a sequence of simple measurable functions that approximates f. Also fix $x \in \mathcal{D}_f$. Then by the Dominated Convergence Theorem, it follows from (3.5) that $(\Psi(s_n)x)$ is a Cauchy sequence in \mathcal{H} , hence it converges to some $y \in \mathcal{H}$. Now let $(s'_n) \subseteq \mathcal{SM}(\mathcal{E})$ be another sequence that approximates f, and let $y' = \lim_{n \to \infty} \Psi(s'_n)x$. The interlaced sequence $s_1, s'_1, s_2, s'_2, \ldots$ also clearly approximates f, so passing to subsequences shows that y = y'. Hence we may define

$$\Psi(f)x = \lim_{n \to \infty} \Psi(s_n)x \tag{3.6}$$

for all $x \in \mathcal{D}_f$. Linearity of $\Psi(f)$ then follows from linearity along an approximating sequence. The general case of the identities (3.3) and (3.4) follows by the Dominated Convergence Theorem and continuity of the inner product and norm.

Finally we note that $\Psi(f)$ is uniquely determined by (3.3) and its domain by polarisation.

REMARK 3.3. For any subset $A \subseteq \Omega$ we may consider the measurable space (A, \mathcal{E}_A) , where \mathcal{E}_A is the σ -algebra inherited from \mathcal{E} . Given a spectral measure $E \colon \mathcal{E}_A \to \mathcal{H}$, we may extend this to \mathcal{E} by letting $\tilde{E}(B) = E(A \cap B)$ for $B \in \mathcal{E}$. If further $A \in \mathcal{E}$, we have for any $f \in \mathcal{M}(\mathcal{E})$,

$$\int_{\Omega} f \, d\tilde{E} = \int_{A} f \, d\tilde{E} + \int_{\Omega \setminus A} f \, d\tilde{E} = \int_{A} f \, dE.$$

PROPOSITION 3.4: Change of variables in spectral integrals

Let (Ω, \mathcal{E}) and (Ω', \mathcal{E}') be measure spaces, and let E be a spectral measure on \mathcal{E} . Let further $\varphi \colon \Omega \to \Omega'$ is $(\mathcal{E}, \mathcal{E}')$ -measurable and let $E' = \varphi_* E$ denote the pushforward of E by φ . Then for all $x \in \mathcal{H}$ and $f \in \mathcal{M}(\mathcal{E}')$, $f \in L^2(E'_x)$ if and only if the pullback

 $\varphi^* f = f \circ \varphi \in L^2(E_x)$. Furthermore,

$$\int f \, \mathrm{d}(\varphi_* E) = \int \varphi^* f \, \mathrm{d}E. \tag{3.7}$$

PROOF. First notice that E_x' is equal to the pushforward $\varphi_*(E_x)$ for all $x \in \mathcal{H}$. If $f \in \mathcal{M}(\mathcal{E}')$, we then have $f \in L^1(E_x')$ if and only if $f \circ \varphi \in L^1(E_x)$ by Thorbjørnsen (2014, Sætning 11.1.4). Furthermore, $f \in L^2(E_x')$ if and only if $|f|^2 \in L^1(E_x')$, which is equivalent to $|f \circ \varphi|^2 = |f|^2 \circ \varphi \in L^1(E_x)$. But that is the case if and only if $f \circ \varphi \in L^2(E_x)$. This proves the first statement. In particular,

$$\left\{x \in \mathcal{H} \left| \int |f|^2 dE' < \infty \right\} = \left\{x \in \mathcal{H} \left| \int |f \circ \varphi|^2 dE < \infty \right\},\right.$$

so each of the two spectral integrals in (3.7) have the same domain.

Furthermore, it also follows from Thorbjørnsen (2014, Sætning 11.1.4) that

$$\int f \, \mathrm{d}E_x' = \int f \circ \varphi \, \mathrm{d}E_x$$

for all x that lie in the common domain of the two integrals. Finally, (3.7) follows by polarisation.

PROPOSITION 3.5: Additivity of spectral integrals

Let $f,g \in \mathcal{M}(\mathcal{E})$. Then

$$\Psi(f) + \Psi(g) \subseteq \Psi(f + g)$$
.

PROOF. Let $x \in \mathcal{D}_f \cap \mathcal{D}_g$. Pick approximating sequences (s_n) , $(t_n) \subseteq \mathcal{SM}(\mathcal{E})$ for f and g respectively and notice that

$$\Psi(s_n + t_n)x = \Psi(s_n)x + \Psi(t_n)x$$
, $(n \in \mathbb{N})$.

We clearly have $\mathcal{D}_f \cap \mathcal{D}_g \subseteq \mathcal{D}_{f+g}$ by a straightforward application of (3.2). So $x \in \mathcal{D}_{f+g}$, and $(s_n + t_n)$ is a sequence of simple functions that approximates f + g, so the desired equality follows by (3.6).

THEOREM 3.6: The multiplication theorem

Let $f,g \in \mathcal{M}(\mathcal{E})$. Then

$$\Psi(f)\Psi(g) \subseteq \Psi(fg)$$
 and $\mathcal{D}(\Psi(f)\Psi(g)) = \mathcal{D}_g \cap \mathcal{D}_{fg}$,

with equality if and only if $\mathcal{D}_{fg} \subseteq \mathcal{D}_g$.

PROOF. Let $x \in \mathcal{D}(\Psi(f)\Psi(g))$. Then $x \in \mathcal{D}_g$ and $y = \Psi(g)x \in \mathcal{D}_f$. Let $(s_n), (\tilde{t}_k) \subseteq \mathcal{SM}(\mathcal{E})$ be sequences that approximate f and g respectively. Now fix $n \in \mathbb{N}$ temporarily. Since $\Psi(s_n)$ is bounded, there exists a $K \in \mathbb{N}$ such that

$$k \ge K \quad \Rightarrow \quad \Psi(s_n)(y - \Psi(\tilde{t}_k)x) < \frac{1}{n}.$$

Choose a subsequence $(t_n) = (\tilde{t}_{k_n})$ that satisfies this property. Then

$$\Psi(f)y = \lim_{n \to \infty} \Psi(s_n)y = \lim_{n \to \infty} \Psi(s_n)\Psi(t_n)x = \lim_{n \to \infty} \Psi(s_nt_n)x.$$
 (3.8)

Fatou's lemma then yields

$$\int |fg|^2 dE_x \le \liminf_{n \to \infty} \int |s_n t_n|^2 dE_x = \liminf_{n \to \infty} ||\Psi(s_n t_n)x||^2 = ||\Psi(f)y||^2 < \infty,$$

since $y \in \mathcal{D}_f$, which shows that $x \in \mathcal{D}_{fg}$. Thus $\mathcal{D}(\Psi(f)\Psi(g)) \subseteq \mathcal{D}_g \cap \mathcal{D}_{fg}$. Continuing (3.8) we get

$$\Psi(f)y = \lim_{n \to \infty} \Psi(s_n t_n) x = \Psi(fg) x,$$

showing that $\Psi(f)\Psi(g) \subseteq \Psi(fg)$.

It remains to show that $\mathcal{D}_g \cap \mathcal{D}_{fg} \subseteq \mathcal{D}(\Psi(f)\Psi(g))$, so let $x \in \mathcal{D}_g \cap \mathcal{D}_{fg}$. Letting $y = \Psi(g)x$ as before, we need to show that $y \in \mathcal{D}_f$. So let $(s_n) \subseteq \mathcal{SM}(\mathcal{E})$ be a sequence that approximates f. Then $\Psi(s_n)y = \Psi(s_ng)x$, so

$$\int |s_n|^2 dE_y = ||\Psi(s_n)y||^2 = ||\Psi(s_ng)x||^2 \le \int |fg|^2 dE_x < \infty.$$

It then follows from Fatou's lemma that $y \in \mathcal{D}_f$ as desired.

The final claim follows easily by considering each direction separately. □

THEOREM 3.7

Let $f \in \mathcal{M}(\mathcal{E})$. Then

$$\Psi(f)^* = \Psi(\overline{f}),\tag{3.9}$$

and

$$\Psi(f)\Psi(f)^* = \Psi(|f|^2) = \Psi(f)^*\Psi(f). \tag{3.10}$$

In particular, $\Psi(f)$ is closed and normal.

PROOF. First let $(s_n) \subseteq \mathcal{SM}(\mathcal{E})$ be a sequence that approximates f. Then for $x, y \in \mathcal{D}_f = \mathcal{D}_{\overline{f}}$ it follows from Theorem C.3 that

$$\langle \Psi(f)x,y\rangle = \lim_{n\to\infty} \langle \Psi(s_n)x,y\rangle = \lim_{n\to\infty} \langle x,\Psi(\overline{s_n})y\rangle = \langle x,\Psi(\overline{f})y\rangle,$$

showing that $\Psi(\overline{f}) \subseteq \Psi(f)^*$.

Now let $y \in \mathcal{D}(\Psi(f)^*)$. Let $A_n = \{|f| \le n\}$ as in the proof of Lemma 3.1. From the proof of that lemma we know that $\mathcal{R}(\Psi(\mathbf{1}_{A_n})) = \mathcal{R}(E(A_n)) \subseteq \mathcal{D}_f$. And since $\mathbf{1}_{A_n}$ is bounded, we have $\Psi(f)\Psi(\mathbf{1}_{A_n}) = \Psi(f\mathbf{1}_{A_n})$ by Theorem 3.6. For all $x \in \mathcal{H}$ we then have

$$\langle x, \Psi(\mathbf{1}_{A_n})\Psi(f)^*y\rangle = \langle \Psi(f)\Psi(\mathbf{1}_{A_n})x, y\rangle = \langle \Psi(f\mathbf{1}_{A_n})x, y\rangle = \langle x, \Psi(\overline{f}\mathbf{1}_{A_n})y\rangle,$$

so $\Psi(\mathbf{1}_{A_n})\Psi(f)^*y = \Psi(\overline{f}\mathbf{1}_{A_n})y$. And by Theorem C.3, $\|\Psi(\mathbf{1}_{A_n})\| = \|\mathbf{1}_{A_n}\|_{\infty} \le 1$, so

$$\int |f|^2 \mathbf{1}_{A_n} dE_y = \|\Psi(\overline{f} \mathbf{1}_{A_n})y\|^2 \le \|\Psi(\mathbf{1}_{A_n})\|^2 \|\Psi(f)^*y\|^2 \le \|\Psi(f)^*y\|^2.$$

Applying the Monotone Convergence Theorem then yields that $y \in \mathcal{D}_{\overline{f}}$, so (3.9) is proved.

The statement (3.10) follows by (3.9) and Theorem 3.6, since $\mathcal{D}_{f\overline{f}}\subseteq \mathcal{D}_f.\square$

PROPOSITION 3.8

Let $f \in \mathcal{M}(\mathcal{E})$. Then $\Psi(f)$ is bounded if and only if $f \in L^{\infty}(E)$.

PROOF. It is clear that if $f \in L^{\infty}(E)$, then $\Psi(f)$ is bounded. Conversely, if $\Psi(f)$ is bounded, then let $A_n = \{|f| \ge ||\Psi(f)|| + 1/n\}$ for $n \in \mathbb{N}$. Then for $x \in \mathcal{H}$ we have

$$\|\Psi(f)\|^{2}\|E(A_{n})x\|^{2} \ge \|\Psi(f)E(A_{n})x\|^{2} = \|\Psi(f\mathbf{1}_{A_{n}})x\|^{2} = \int_{A_{n}} |f|^{2} dE_{x}$$

$$\ge \left(\|\Psi(f)\| + \frac{1}{n}\right)^{2} \|E(A_{n})x\|^{2},$$

which is only possible if $E(A_n) = 0$. Letting $A = \bigcup_{n \in \mathbb{N}} A_n = \{|f| > ||\Psi(f)||\}$ we then have E(A) = 0, so $|f| \le ||\Psi(f)||$ E-a.e., proving the claim.

PROPOSITION 3.9

Let $f \in \mathcal{M}(\mathcal{E})$. Then $\Psi(f)$ is invertible if and only if $f \neq 0$ E-a.e. In that case $\Psi(f)^{-1} = \Psi(g)$ for $g = (1/f)\mathbf{1}_{\Omega \setminus N}$, where $N = \{f = 0\}$.

PROOF. First notice that $\Psi(f)E(N) = \Psi(f\mathbf{1}_N) = 0$, so if $E(N) \neq 0$, then $\Psi(f)$ is not invertible.

Conversely, assume that E(N)=0. Then $\Psi(gf)=\Psi(1)=I$, so $\mathcal{D}_{gf}=\mathcal{H}$. It follows from Theorem 3.6 that $\mathcal{D}(\Psi(g)\Psi(f))=\mathcal{D}_f$ with $\Psi(g)\Psi(f)\subseteq\Psi(gf)=I$, so $\Psi(f)$ is injective on its domain. Hence it is invertible with $\Psi(f)^{-1}\subseteq\Psi(g)$. The opposite inclusion follows by symmetry, proving the claim.

For a function $f \in \mathcal{M}(\mathcal{E})$, the essential range of f with respect to E is the set

$$\operatorname{sp}_E f = \{ \lambda \in \mathbb{C} \mid \forall \varepsilon > 0 \colon E(\{|f - \lambda| < \varepsilon\}) \neq 0 \}.$$

It turns out that this can be used to characterise the spectrum of spectral integrals:

PROPOSITION 3.10

Let $f \in \mathcal{M}(\mathcal{E})$. Then we have $\sigma(\Psi(f)) = \operatorname{sp}_E f$.

PROOF. By definition $0 \in \rho(\Psi(f))$ if and only if $\Psi(f)$ is boundedly invertible. But by the previous two propositions, this is the case if and only if E(N) = 0 for $N = \{f = 0\}$, and $g = (1/f)\mathbf{1}_{\Omega\setminus N} \in L^{\infty}(E)$. This is equivalent to the existence of a t > 0 such that $E(\{|g| > t\}) = 0$, i.e., the existence of an $\varepsilon > 0$ such that $E(\{|f| < \varepsilon\}) = 0$, by Skibsted (2019, Lemma 3.45).

The case for general $\lambda \in \mathbb{C}$ follows by replacing f by $f - \lambda$.

3.2 • The spectral theorem for self-adjoint operators

THEOREM 3.11: The spectral theorem for self-adjoint operators

For every self-adjoint operator H in H there exists a unique spectral measure E on the Borel σ -algebra $\mathfrak{B}(\mathbb{R})$ of the real line such that

$$H = \int \mathrm{id}_{\mathbb{R}} \, \mathrm{d}E.$$

The spectral measure E is called the spectral resolution *of H*.

PROOF. We first prove existence of E. Let U be the Cayley transform of E. Then by Theorem C.4 there is a spectral measure E' on the spectrum $\sigma(U)$ of U such that $U = \int \operatorname{id}_{\sigma(U)} \operatorname{d}E'$. Since U is unitary, its spectrum lies on the unit circle $S^1 \subseteq \mathbb{C}$, and so by Remark 3.3 we may extend E' to S^1 . But since I - U is injective by Theorem 2.2, 1 is not an eigenvalue of U, and so $E'(\{1\}) = 0$ by Skibsted (2019, Lemma 3.70) or Rudin (1991, Theorem 12.29(b)). Letting $\Omega = S^1 \setminus \{1\}$, it then follows that

$$\int_{\sigma(U)} f \, \mathrm{d}E' = \int_{\Omega} f \, \mathrm{d}E'$$

for all $f \in \mathcal{M}(\mathfrak{B}(\Omega))$. In particular, $U = \int \mathrm{id}_{\Omega} \, \mathrm{d}E'$. We thus consider E' as a spectral measure on $\mathfrak{B}(\Omega)$.

Now consider the Möbius transformation $\psi \colon \Omega \to \mathbb{R}$ given by $\psi(\omega) = \mathrm{i}(1 + \omega)(1-\omega)^{-1}$. This is clearly a homeomorphism with inverse $\varphi(\lambda) = (\lambda - \mathrm{i})(\lambda + \mathrm{i})^{-1}$.

Then define a spectral measure $E = \psi_* E'$ on $\mathfrak{B}(\mathbb{R})$ as the pushforward of E' by ψ . Define an operator $\tilde{H} = \int \mathrm{id}_{\mathbb{R}} dE$. By Proposition 3.4 we have

$$\int f \, \mathrm{d}E = \int f \circ \psi \, \mathrm{d}E'$$

for all $f \in \mathcal{M}(\mathfrak{B}(\mathbb{R}))$. Consider in particular the function $f(\lambda) = (\lambda - i)^{-1}$. Then $(f \circ \psi)(\omega) = i(1 - \omega^{-1})/2$, which is clearly bounded. It then follows that

$$(\tilde{H} - i)^{-1} = \int f(\lambda) dE(\lambda) = \frac{i}{2} \int (1 - \omega^{-1}) dE'(\omega)$$
$$= \frac{i}{2} (1 - U^{-1}) = \frac{i}{2} (1 - (H + i)(H - i)^{-1}) = (H - i)^{-1}$$

Hence $H = \tilde{H}$, and so $H = \int id_{\mathbb{R}} dE$ as desired.

We now turn to uniqueness. Let E be a spectral measure on $\mathfrak{B}(\mathbb{R})$ such that $H = \int \mathrm{id}_{\mathbb{R}} \, \mathrm{d}E$, and define a spectral measure $E' = \varphi_* E$ on $\mathfrak{B}(\Omega)$ as before. Then since id_{Ω} is bounded, the operator $U = \int \mathrm{id}_{\Omega} \, \mathrm{d}E'$ is also bounded. But then

$$U = \int \mathrm{id}_{\Omega} \, \mathrm{d}E' = \int \varphi \, \mathrm{d}E = (H - \mathrm{i})(H + \mathrm{i})^{-1},$$

So U is uniquely determined by H. But the uniqueness part of Theorem C.4 then yields that E' is unique given H. Finally, E is uniquely determined by E' since $E = \psi_* E'$.

PROPOSITION 3.12: Spectral resolution and unitary equivalence

Let T be a self-adjoint operator on \mathcal{H} with spectral resolution E^T , U a unitary operator, and define $S = U^*TU$. Then S is self-adjoint with spectral resolution E^S given by $E^S(A) = U^*E^T(A)U$ for $A \in \mathcal{E}$, and $E_x^S = E_{Ux}^T$ for all $x \in \mathcal{H}$. Furthermore, for all $f \in \mathcal{M}(\mathcal{E})$ we have

$$\int f \, \mathrm{d}E^S = U^* \int f \, \mathrm{d}E^T \, U. \tag{3.11}$$

If two operators *T* and *S* are related as above, we say that they are *unitarily equivalent*.

PROOF OF PROPOSITION 3.12. It is obvious that S is self-adjoint. It is also easy to check that E^S as above defined is a spectral measure, so we omit the proof. We now show that it is the spectral resolution of S. For $x \in \mathcal{H}$ and $A \in \mathcal{E}$ we have

$$E_x^S(A) = \langle E^S(A)x, x \rangle = \langle E^S(A)Ux, Ux \rangle = E_{Ux}^T(A)$$
 (3.12)

as claimed. It is then clear that the domain of the operator $\int id_{\mathbb{R}} dE^S$ is exactly $\mathcal{D}(S)$. Thus we get, for $x \in \mathcal{D}(S)$,

$$\left\langle \int i d_{\mathbb{R}} dE^{S} x, x \right\rangle = \int i d_{\mathbb{R}} dE_{x}^{S} = \int i d_{\mathbb{R}} dE_{Ux}^{T} = \left\langle \int i d_{\mathbb{R}} dE^{T} Ux, Ux \right\rangle$$
$$= \left\langle TUx, Ux \right\rangle = \left\langle U^{*}TUx, x \right\rangle = \left\langle Sx, x \right\rangle.$$

So $\int id_{\mathbb{R}} dE^S = S$ by polarisation, and E^S is then the spectral resolution of S by the uniqueness part of Theorem 3.11.

Finally we consider the identity (3.11). By (3.12) it is clear that the two operators have the same domains. Equality then follows by a similar calculation to the one above.

_ 4 _

One-parameter unitary groups

One-parameter groups of operators have many applications throughout mathematics and physics. We begin by proving some basic properties of these, after which we turn to one of the main theorems, Stone's Theorem 4.4. We end this section by showing how Stone's theorem can shed some light on the Schrödinger equation.

4.1 • Unitary groups and Stone's theorem

DEFINITION 4.1: Unitary groups

A one-parameter unitary group is a group homomorphism $U: (\mathbb{R}, +) \to U(\mathcal{H})$. We say that U is *strongly continuous* if for all $t \in \mathbb{R}$ and $x \in \mathcal{H}$ we have

$$\lim_{h\to 0} U(t+h)x = U(t)x.$$

Spelling out the homomorphism property, this means that U(s+t) = U(s)U(t) for all $s, t \in \mathbb{R}$. In particular, U(0) = I and $U(-t) = U(t)^{-1} = U(t)^*$. The strong continuity will usually be implicit, so we will often simply use the term 'unitary group'.

REMARK 4.2. It turns out that strong continuity of one-parameter unitary groups is equivalent to several seemingly weaker properties. Firstly it is clear that it is enough to assume strong continuity for t = 0 and for x in a dense subset of \mathcal{H} . A short calculation shows that it is also enough to assume weak continuity. If furthermore \mathcal{H} is separable, it is even enough for U to be weakly Lebesgue-measurable, i.e. for the map $t \mapsto \langle U(t)x,y \rangle$ to be Lebesgue-measurable for all $x,y \in \mathcal{H}$. This is a theorem due to von Neumann, see Reed and Simon (1980, Theorem VIII.9).

If H is a self-adjoint operator on \mathcal{H} with spectral resolution E, the spectral theorem allows us to define the operator

$$e^{itH} = \int_{\mathbb{R}} e^{it\lambda} dE(\lambda)$$

for $t \in \mathbb{R}$. As expected this family of operators constitutes a unitary group:

PROPOSITION 4.3: Properties of unitary groups

Let H be a self-adjoint operator on H, and let $U: \mathbb{R} \to \mathcal{H}$ be the map $U(t) = e^{itH}$. Then the following statements are true:

- (i) The map U is a strongly continuous one-parameter unitary group.
- (ii) The operator H is uniquely determined by U in the following sense: The domain $\mathcal{D}(H)$ consists exactly of those $x \in \mathcal{H}$ such that the derivative

$$\frac{\mathrm{d}}{\mathrm{d}t}\bigg|_{t=0} U(t)x = \lim_{h \to 0} \frac{(U(h) - I)x}{h} \tag{4.1}$$

exists, and the derivative equals i Hx if it exists.

(iii) If $x \in \mathcal{D}(H)$ and $t \in \mathbb{R}$, then $U(t)x \in \mathcal{D}(H)$ and

$$\frac{\mathrm{d}}{\mathrm{d}t}U(t)x = \mathrm{i}\,HU(t)x = \mathrm{i}\,U(t)Hx. \tag{4.2}$$

PROOF. It is clear from the spectral theorem that U is a group homomorphism into $U(\mathcal{H})$. To see that it is strongly continuous, notice that, for $x \in \mathcal{H}$,

$$||U(h)x - x||^2 = \int_{\mathbb{R}} |e^{ih\lambda} - 1|^2 dE_x(\lambda) \to 0$$
 as $h \to 0$

by the Dominated Convergence Theorem, since the integrand is bounded and the measure E_x is finite. This proves (i).

To prove (ii) we first let $x \in \mathcal{D}(H)$. If we write the function $\lambda \mapsto e^{ih\lambda} - 1$ as the integral of its derivative, we find that $|e^{ih\lambda} - 1| \le |h\lambda|$. Thus

$$\left\|\frac{(U(h)-I)x}{h}-\mathrm{i}\,Hx\right\|^2=\int_{\mathbb{R}}\left|\frac{\mathrm{e}^{\mathrm{i}\,h\lambda}-1}{h}-\mathrm{i}\,\lambda\right|^2\mathrm{d}E_x(\lambda)\leq 4\int_{\mathbb{R}}\lambda^2\,\mathrm{d}E_x(\lambda)<\infty,$$

where the final inequality follows since $x \in \mathcal{D}(H)$. The Dominated Convergence Theorem then implies that the first integral tends to zero as h goes to zero. This proves (4.1) with the correct limit.

Conversely, define an operator T on \mathcal{H} by

$$Tx = \lim_{h \to 0} \frac{1}{i} \frac{(U(h) - I)x}{h}, \quad (x \in \mathcal{D}(T)),$$
 (4.3)

where the domain $\mathcal{D}(T)$ is the set of those $x \in \mathcal{H}$ for which the limit (4.1) exists. The above then shows that $H \subseteq T$. Let $x, y \in \mathcal{D}(T)$ and notice that

$$\langle Tx, y \rangle = \lim_{h \to 0} \left\langle \frac{1}{i} \frac{(U(h) - I)x}{h}, y \right\rangle = \lim_{h \to 0} \left\langle x, \frac{1}{i} \frac{(U(-h) - I)y}{-h} \right\rangle = \langle x, Ty \rangle,$$

so T is symmetric. But since H is maximally symmetric, we have H = T. In total, this proves (ii).

Finally we prove (iii). Let $x \in \mathcal{H}$ and $t \in \mathbb{R}$ and notice that

$$\lim_{h \to 0} \frac{(U(h) - I)U(t)x}{h} = U(t)\lim_{h \to 0} \frac{(U(h) - I)x}{h} = U(t)iHx,$$

so $U(t)x \in \mathcal{D}(H)$. The left-hand side equals iHU(t)x by (ii), proving the claim.

THEOREM 4.4: Stone's theorem

If U is a strongly continuous one-parameter unitary group on \mathcal{H} , then there exists a unique self-adjoint operator H on \mathcal{H} such that $U(t) = e^{itH}$ for $t \in \mathbb{R}$.

The operator H is called the *infinitesimal generator* of U.⁶

This theorem was first announced in Stone (1930) and was proved in Stone (1932b). In the proof of Stone's theorem, some authors use integrals of Hilbert space-valued functions⁷, see e.g. Reed and Simon (1980, Theorem VIII.8), Conway (1990, Theorem X.5.6) and Schmüdgen (2012, Theorem 6.2). We take

 $^{^6}$ Some authors, see e.g. Schmüdgen (2012, Theorem 6.2), use the term 'infinitesimal generator' to refer to the operator iH.

 $^{^{7}}$ Since we only wish to integrate continuous \mathcal{H} -valued functions with compact support, it is possible to simply define a Riemann integral. In the present presentation we prefer the following argument using bounded sesquilinear forms, since this avoids proving or postulating the existence of such an integral, and it is not significantly longer.

a slightly different approach following Hall (2013, Lemma 10.18), which is in fact closer to Stone's original proof: If U is a unitary group on \mathcal{H} and $f \in L^1(\mathbb{R})$ is real-valued, then the map $\Gamma_f \colon \mathcal{H} \times \mathcal{H} \to \mathbb{C}$ given by

$$\Gamma_f(x,y) = \int_{\mathbb{R}} f(t) \langle U(t)x,y \rangle dt, \quad (x,y \in \mathcal{H}),$$

is clearly a bounded sesquilinear form. By Theorem B.3 there exists a (unique) bounded operator B_f on \mathcal{H} such that $\Gamma_f(x,y) = \langle B_f x, y \rangle$.

PROOF OF THEOREM 4.4. First define a symmetric operator T as in (4.3). We wish to show that T is densely defined. The first step is to show that for all $f \in C_c^{\infty}(\mathbb{R})$, we have $\mathcal{R}(B_f) \subseteq \mathcal{D}(T)$. For $x, y \in \mathcal{H}$ and $h \in \mathbb{R}$ we have

$$\langle (U(h) - I)B_f x, y \rangle = \langle B_f x, (U(h)^* - I)y \rangle = \int_{\mathbb{R}} f(t) \langle U(t)x, (U(h)^* - I)y \rangle dt$$

$$= \int_{\mathbb{R}} f(t) \langle (U(t+h) - U(t))x, y \rangle dt$$

$$= \int_{\mathbb{R}} (f(t-h) - f(t)) \langle U(t)x, y \rangle dt,$$

where the final integral arises from making a change of variable in the first term of the previous integral. The integrand is continuous and has compact support, so the Dominated Convergence Theorem yields

$$\lim_{h\to 0} \left\langle \frac{(U(h)-I)B_f x}{h}, y \right\rangle = -\int_{\mathbb{R}} f'(t) \langle U(t)x, y \rangle dt = -\langle B_{f'}x, y \rangle,$$

showing that $B_f x \in \mathcal{D}(T)$ as desired.

Next, choose a sequence $(f_n) \subseteq C_c^{\infty}(\mathbb{R})$ of non-negative functions supported in $I_n = [-1/n, 1/n]$ and with the property that $\int_{\mathbb{R}} f_n(t) dt = 1$. A short calculation shows that, for $x, y \in \mathcal{H}$,

$$\langle B_{f_n} x - x, y \rangle = \int_{\mathbb{R}} f_n(t) \langle (U(t) - I)x, y \rangle dt.$$

Letting $y = B_{f_n}x - x$, it follows by the Cauchy–Schwarz inequality that

$$||B_{f_n}x - x|| \le \int_{\mathbb{R}} f_n(t) ||(U(t) - I)x|| dt \le \sup_{t \in I_n} ||U(t)x - x||,$$

which converges to zero as t tends to zero since U is strongly continuous. Thus $\mathcal{D}(T)$ is dense in \mathcal{H} .

We now show that T is essentially self-adjoint. Let $y \in \mathcal{N}(T^* - i)$. From the definition of T, it follows that $U(t)\mathcal{D}(T) \subseteq \mathcal{D}(T)$, so letting $x \in \mathcal{D}(T)$ we get

$$\frac{\mathrm{d}}{\mathrm{d}t}\bigg|_{t=t_0} \langle U(t)x, y \rangle = \frac{\mathrm{d}}{\mathrm{d}s}\bigg|_{s=0} \langle U(s)U(t_0)x, y \rangle = \langle \mathrm{i}\,T\,U(t_0)x, y \rangle$$
$$= \langle \mathrm{i}\,U(t_0)x, T^*y \rangle = \langle U(t_0)x, y \rangle.$$

So the function $g(t) = \langle U(t)x, y \rangle$ solves the differential equation g' = g on \mathbb{R} , and thus $g(t) = g(0)e^t$. But g is bounded, so we must have g(0) = 0, and hence $\langle x, y \rangle = 0$. Since $\mathcal{D}(T)$ is dense, we have y = 0, and so $\mathcal{N}(T^* - i) = 0$. It similarly follows that $\mathcal{N}(T^* + i) = 0$, and Proposition 1.16 then shows that T is essentially self-adjoint.

Define a unitary group V by $V(t) = e^{i t \overline{T}}$. We aim to show that U = V. For $x \in \mathcal{D}(T)$ we get

$$\frac{\mathrm{d}}{\mathrm{d}t}\bigg|_{t=t_0}U(t)x = \frac{\mathrm{d}}{\mathrm{d}s}\bigg|_{s=0}U(s)U(t_0)x = \mathrm{i}\,T\,U(t_0)x,$$

and similarly, $\frac{d}{dt}|_{t=t_0}V(t)x=i\overline{T}V(t_0)x$. Now let w(t)=(U(t)-V(t))x. Then $\frac{d}{dt}w(t)=i\overline{T}w(t)$. It follows that

$$\frac{\mathrm{d}}{\mathrm{d}t}\|w(t)\|^2 = \frac{\mathrm{d}}{\mathrm{d}t}\langle w(t), w(t)\rangle = \langle \mathrm{i}\,\overline{T}w(t), w(t)\rangle + \langle w(t), \mathrm{i}\,\overline{T}w(t)\rangle = 0,$$

since \overline{T} is symmetric. But w(0) = 0, so w(t) = 0 for all $t \in \mathbb{R}$. It follows that U = V on $\mathcal{D}(T)$, and since $\mathcal{D}(T)$ is dense in \mathcal{H} , U = V on \mathcal{H} . Letting $H = \overline{T}$, the claim follows.

4.2 • The Schrödinger equation

Let $I \subseteq \mathbb{R}$ be an open interval, and consider a Hilbert space-valued function $\psi \colon I \to \mathcal{H}$. If the limit $\lim_{h\to 0} h^{-1}(\psi(t+h)-\psi(t))$ exists for some $t\in I$, we say that ψ is differentiable at t and denote the limit $\psi'(t)$. Then let $C^1(I,\mathcal{H})$ be the set of all functions $\psi \colon I \to \mathcal{H}$ that are differentiable at every point in I such that the function ψ' is continuous.

Let H be a self-adjoint operator on \mathcal{H} . Consider the abstract Schrödinger equation⁸

$$\psi'(t) = -iH\psi(t), \quad (t \in \mathbb{R}). \tag{4.4}$$

$$i\hbar \frac{\mathrm{d}}{\mathrm{d}t}|\psi(t)\rangle = \widehat{H}|\psi(t)\rangle,$$

for some concrete Hamiltonian operator \widehat{H} .

 $^{^8}$ Compare this to the time-dependent Schrödinger equation as is usually appears in physics:

Given an initial condition $\psi_0 \in \mathcal{H}$, we wish to solve the *Cauchy problem* for (4.4): Find a $\psi \in C^1(\mathbb{R}, \mathcal{H})$ such that $\psi(t) \in \mathcal{D}(H)$ for all $t \in \mathbb{R}$, ψ satisfies (4.4), and $\psi(0) = \psi_0$.

PROPOSITION 4.5: The abstract Schrödinger equation

For any $\psi_0 \in \mathcal{D}(H)$ there is a unique solution to the Cauchy problem for the abstract Schrödinger equation given by

$$\psi(t) = e^{-itH} \psi_0, \quad (t \in \mathbb{R}).$$

PROOF. That ψ is a solution follows from Proposition 4.3(iii). As for uniqueness, let ψ_1 and ψ_2 be two solutions and let $\varphi = \psi_1 - \psi_2$. Then $\varphi'(t) = -\mathrm{i}\,H\varphi(t)$ and $\varphi(0) = 0$. Thus

$$\frac{\mathrm{d}}{\mathrm{d}t} \|\varphi(t)\|^2 = \langle \varphi'(t), \varphi(t) \rangle + \langle \varphi(t), \varphi'(t) \rangle$$
$$= \langle -\mathrm{i} H \varphi(t), \varphi(t) \rangle + \langle \varphi(t), -\mathrm{i} H \varphi(t) \rangle = 0,$$

where the final equality follows since H is symmetric. But then $\varphi(t) = 0$ for all $t \in \mathbb{R}$, so $\psi_1 = \psi_2$.

Notice also that this confirms the well-known result from quantum mechanics that the dynamics given by the Schrödinger equation is unitary. So the Schrödinger equation is consistent with the usual requirement that wave functions are normalised, and thus consistent with the statistical interpretation.

Conversely, if one takes the statistical interpretation as fundamental, and thus requires the dynamics to be unitary (among other reasonable physical assumptions), then the time evolution of a wave function is given by a unitary group U. That is, for an initial condition $\psi_0 \in \mathcal{H}$ we may define the continuous map $\psi(t) = U(t)\psi_0$. Stone's theorem then guarantees that the dynamics is generated by some self-adjoint operator H, the *Hamiltonian*. We thus recover the abstract Schrödinger equation.

We refer to Sakurai and Napolitano (2011, Chapter 2) for further discussion of the underlying physics.

— 5 **—**

The operators of quantum mechanics

We now consider some of the concrete operators that appear in quantum mechanics: The position, potential, momentum and kinetic energy operators. The first two are multiplication operators, and we begin by motivating this fact. The latter two are linear differential operators, and we recover these by applying the Fourier transform to appropriate multiplication operators. Finally we take a closer look at the Laplace-operator $-\Delta$ which gives the dynamics of a free particle.

5.1 • Multiplication operators: Position and potential energy

Consider a single particle on the real line. Its (normalised) wave function is a map⁹ ψ : $\mathbb{R} \to \mathbb{C}$, and we interpret the function $|\psi|^2$ as the probability density for the position of the particle. Thus the expectation value of the position x in the state ψ is

$$\mathbb{E}_{\psi}(x) = \int_{\mathbb{R}} x |\psi(x)|^2 \, \mathrm{d}x.$$

If we introduce the position operator X given by $X\psi(x)=x\psi(x)$, we can write the expectation value $\mathbb{E}_{\psi}(x)=\langle X\psi,\psi\rangle$ as we are want to do in quantum mechanics. Similarly, if V is a potential, i.e., a (sufficiently nice) function of position, the expectation value of V is

$$\mathbb{E}_{\psi}(V) = \int_{\mathbb{R}} V(x) |\psi(x)|^2 \, \mathrm{d}x = \langle M_V \psi, \psi \rangle,$$

where M_V is the operator given by $M_V \psi = V \psi$.

Thus we are led to consider multiplication operators, and since these are interesting in their own right, we cast the following theorem in terms of a general multiplication operator:

⁹ This initial discussion is purely heuristic, but already now we think of ψ as living in $L^2(\mathbb{R})$. Also, the space of states of a quantum system is more properly thought of as a *projective Hilbert space* $P(\mathcal{H})$ of one-dimensional (complex) subspaces of a Hilbert space \mathcal{H} . Usually we think of states given by a unit vector in \mathcal{H} , two such vectors describing the same state if they are scalar multiples of each other.

THEOREM 5.1: Multiplication operators

Let $(\Omega, \mathcal{E}, \mu)$ be a measure space, and let $\varphi \in \mathcal{M}(\mathcal{E})$. Define an operator M_{φ} on $L^2(\mu)$ with domain

$$\mathcal{D}(M_{\varphi}) = \{ f \in L^2(\mu) \mid \varphi f \in L^2(\mu) \}$$

by $M_{\varphi}f=\varphi f$. Then M_{φ} is densely defined, $(M_{\varphi})^*=M_{\overline{\varphi}}$ and M_{φ} is normal. If φ is real-valued, then M_{φ} is self-adjoint.

PROOF. Define a spectral measure $E: \mathcal{E} \to \mathcal{B}(L^2(\mu))$ by letting $E(A)f = \mathbf{1}_A f$ for $A \in \mathcal{E}$ and $f \in L^2(\mu)$. If $\varphi \in \mathcal{SM}(\mathcal{E})$, then it is easy to see that $\Psi(\varphi)f = M_{\varphi}f$ for all $f \in L^2(\mu)$. Furthermore,

$$\int_{\Omega} |\varphi|^2 dE_f = ||\Psi(\varphi)f||^2 = \int_{\Omega} |\varphi f|^2 d\mu.$$
 (5.1)

If φ is a general function from $\mathcal{M}(\mathcal{E})$, then choose a sequence $(\varphi_n) \subseteq \mathcal{SM}(\mathcal{E})$ that approximates φ from below. Then replacing φ by φ_n in (5.1) and applying the Dominated Convergence Theorem shows that $\varphi \in L^2(E_f)$ if and only if $\varphi f \in L^2(\mu)$, so $\mathcal{D}_{\varphi} = \mathcal{D}(M_{\varphi})$. But then M_{φ} is densely defined by Lemma 3.1. Finally it follows from $\Psi(\varphi_n)f = \varphi_n f$ that $\Psi(\varphi)f = \varphi f$ by strong continuity of Ψ .

It then follows from Theorem 3.7 that $(M_{\varphi})^* = M_{\overline{\varphi}}$ and that M_{φ} is normal. If φ is real-valued, then $\varphi = \overline{\varphi}$, and so $M_{\varphi} = (M_{\varphi})^*$.

There exist more elementary proofs of Theorem 5.1 (at least for the σ -finite case) that do not use spectral measures, e.g. Schmüdgen (2012, Example 3.8). But the proof above nicely illustrates a connection between general self-adjoint operators (or spectral measures) and multiplication operators. In fact, any self-adjoint operator on a separable Hilbert space is unitarily equivalent to a multiplication operator, see for example Reed and Simon (1980, Theorem VIII.4) or Hall (2013, Theorem 10.10).

We now characterise the spectrum of multiplication operators. A subset A of a measure space $(\Omega, \mathcal{E}, \mu)$ is called an atom if for any measurable $B \subseteq A$ with $\mu(B) < \mu(A)$, we have $\mu(B) = 0$. An ∞ -atom A is an atom with $\mu(A) = \infty$. We say that μ is semifinite if for any $A \in \mathcal{E}$ with $\mu(A) = \infty$, there exists a measurable $B \subseteq A$ with $0 < \mu(B) < \infty$. Notice that any σ -finite measure is semifinite. Furthermore, the $essential\ range^{11}\ {\rm sp}_{\mu}\ \varphi$ of a function $\varphi \in \mathcal{M}(\mathcal{E})$ with respect to the measure μ is the complement of the union of all open $U \subseteq \mathbb{C}$ such that $\mu(\varphi^{-1}(U)) = 0$. If μ is the Lebesgue measure, we simply write ${\rm sp}\ \varphi$.

¹⁰ In the sequel we will only be concerned with σ-finite measures, so the reader may without issue substitute 'σ-finite' for 'semifinite' in the statement and proof of the following proposition.

¹¹ Not to be confused with the essential range $\operatorname{sp}_E \varphi$ of φ with respect to a spectral measure E.

PROPOSITION 5.2: The spectrum of a multiplication operator

The resolvent $\rho(M_{\varphi})$ is the union of $\Omega \setminus \operatorname{sp}_{\mu} \varphi$ with all open sets $U \subseteq \mathbb{C}$ such that $\varphi^{-1}(U)$ is an ∞ -atom. If μ is semifinite, then $\sigma(M_{\varphi}) = \operatorname{sp}_{\mu} \varphi$.

PROOF. Let \mathcal{U}_{∞} be the set of all such U, and let E be the spectral measure from the proof of Theorem 5.1. We claim that for any $A \in \mathcal{E}$, we have E(A) = 0 if and only if $\mu(A) = 0$ or A is an ∞ -atom. It is clear that if $\mu(A) = 0$, then E(A) = 0. Assume then that A an ∞ -atom, and let $f \in L^2(\mu)$. For $n \in \mathbb{N}$, define the set $A_n = \{|f|^2 \mathbf{1}_A \ge 1/n\} \subseteq A$. Then this set is a μ -null set, since otherwise we would have $\mu(A_n) = \infty$, and then

$$\int_{\Omega} |f|^2 d\mu \ge \int_{A_n} |f|^2 d\mu \ge \frac{1}{n} \mu(A_n) = \infty,$$

contradicting that $f \in L^2(\mu)$. But then the set

$$\{|f|^2 \mathbf{1}_A \neq 0\} = \bigcup_{n \in \mathbb{N}} \{|f|^2 \mathbf{1}_A \ge 1/n\}$$

is also a null set, and so $f = f \mathbf{1}_{\Omega \setminus A} \mu$ -a.e. It follows that $E(A)f = f \mathbf{1}_A = 0$.

Conversely, assume that E(A) = 0 with A not an ∞ -atom. Then either $\mu(A) = 0$ or there is some subset $B \subseteq A$ with $0 < \mu(B) < \infty$. But then $0 \ne \mathbf{1}_B \in L^2(\mu)$, and $0 = \mathbf{1}_A \mathbf{1}_B = \mathbf{1}_B$, which is impossible.

Applying the above to sets $A = \varphi^{-1}(U)$ for open $U \subseteq \mathbb{C}$, we get

$$\Omega \setminus \operatorname{sp}_E \varphi = (\Omega \setminus \operatorname{sp}_{\mu} \varphi) \bigcup_{U \in \mathcal{U}_{\infty}} U,$$

and the first claim follows from Proposition 3.10.

The last claim follows since semifinite spaces have no ∞ -atoms.

EXAMPLE 5.3: Potential operators.

Let $V \colon \mathbb{R}^d \to \mathbb{R}$ be a Lebesgue-measurable function. We think of V as a potential function. Then the corresponding multiplication operator $M_V \psi = V \psi$ is self-adjoint on its domain $\mathcal{D}(M_V) = \{f \in L^2(\mathbb{R}^d) \mid V \psi \in L^2(\mathbb{R}^d)\}$. Since the Lebesgue measure is σ -finite, the spectrum is $\sigma(M_V) = \operatorname{sp} V$.

EXAMPLE 5.4: Position operators.

The position operator X_j given by $X_j\psi(x)=x_j\psi(x)$ is a special case of a potential operator, where V is the jth coordinate projection. It is thus also self-adjoint with spectrum $\sigma(X_j)=\mathbb{R}$.

5.2 • Differential operators: Momentum and kinetic energy

We now turn to differential operators, in particular momentum and kinetic energy operators. We first consider general linear differential operators with constant coefficients. Our presentation is a generalisation of Reed and Simon (1975, Theorem IX.27) and the surrounding discussion. For further properties of these operators, see Schmüdgen (2012, Section 8.2). We refer to Appendix A for an overview of notation.

Let $q(x) = \sum_{|\alpha| \le d} a_{\alpha} x^{\alpha}$ be a polynomial with complex coefficients. Then the corresponding differential operator q(D) acts on functions from $L^2(\mathbb{R}^d)$ as distributions. Define the domains

$$\mathcal{D}_{\max} = \{ \varphi \in L^2(\mathbb{R}^d) \mid q(D)\varphi \in L^2(\mathbb{R}^d) \text{ in the sense of distributions} \},$$

$$\mathcal{D}_{\min} = C_c^{\infty}(\mathbb{R}^d).$$

Let T_{max} be the operator q(D) with domain \mathcal{D}_{max} and T_{min} with domain \mathcal{D}_{min} . Notice that T_{min} is densely defined by Proposition B.1.

Theorem 5.5: Linear differential operators with constant coefficients

- (i) For $\varphi \in L^2(\mathbb{R}^d)$, we have $\varphi \in \mathcal{D}_{max}$ if and only if $q\widehat{\varphi} \in L^2(\mathbb{R}^d)$. We have $T_{max} = \mathcal{F}^{-1}M_q\mathcal{F}$, and T_{max} is normal.
- (ii) The spectrum of T_{max} is its essential range sp q.
- (iii) If q has real coefficients, then T_{max} is self-adjoint.
- (iv) If q has real coefficients, then T_{\min} is essentially self-adjoint with $\overline{T}_{\min} = T_{\max}$.

Notice in particular that T_{max} is unitarily equivalent to the multiplication operator M_a through the Fourier transform.

PROOF OF THEOREM 5.5. Since q is a polynomial, (A.3) shows that $\widehat{q(D)F} = q\widehat{F}$ for any tempered distribution F. Considering φ as a distribution we immediately get the first claim in (i). Furthermore, since \mathcal{F} is unitary and M_q is normal by Theorem 5.1, T_{max} is also normal. Now (ii) follows since $\sigma(T_{\text{max}}) = \sigma(M_q) = \operatorname{sp} q$ by Proposition 5.2, because the spectrum is invariant under unitary equivalence.

From now on assume that q has real coefficients. Then q is real-valued, so T_{max} is self-adjoint by Theorem 5.1, proving (iii).

Finally we show that T_{\min} is essentially self-adjoint. It is sufficient to show that $T_{\min}^* = T_{\max}$, since then $\overline{T}_{\min} = T_{\min}^{**} = T_{\max}$ by Proposition 1.10(iv). First let $\psi \in \mathcal{D}(T_{\min}^*)$. Then we have

$$\langle q(D)\varphi,\psi\rangle = \langle T_{\min}\varphi,\psi\rangle = \langle \varphi,T_{\min}^*\psi\rangle$$

for all $\varphi \in C_c^\infty(\mathbb{R}^d)$, so $q(D)\psi \in L^2(\mathbb{R}^d)$ as a distribution, and then $\psi \in \mathcal{D}_{\max}$ and $T_{\min}^*\psi = q(D)\psi = T_{\max}\psi$. Conversely, let $\psi \in \mathcal{D}_{\max}$. Then $q(D)\psi \in L^2(\mathbb{R}^d)$ as a distribution, and we have $\langle q(D)\varphi,\psi\rangle = \langle \varphi,q(D)\psi\rangle$ for all $\varphi \in C_c^\infty(\mathbb{R}^d)$, showing that $\psi \in \mathcal{D}(T_{\min}^*)$. This proves (iv).

EXAMPLE 5.6: Momentum operators.

There are many ways of motivating the momentum operator. One approach based on Hall (2013, Chapter 3) is to appeal to the *de Broglie hypothesis*, which states (loosely) that a wave function of the form $\psi_k(x) = \mathrm{e}^{\mathrm{i} kx}/\sqrt{2\pi}$ for $k \in \mathbb{Z}$ represents a particle with momentum $p = \hbar k$. This is obviously not square integrable on \mathbb{R} , so for a moment we think of the particle as living on a circle. The functions ψ_k form an orthonormal basis for the Hilbert space $L^2[0,2\pi]$, so the typical wave function for a particle on a circle is on the form

$$\psi(x) = \sum_{k=-\infty}^{\infty} c_k \psi_k(x) = \frac{1}{\sqrt{2\pi}} \sum_{k=-\infty}^{\infty} c_k e^{ikx}.$$

Assuming ψ is normalised, we interpret this as saying that a measurement of the particle's momentum will yield the value $\hbar k$ with probability $|c_k|^2$. The expectation value of the momentum p in the state ψ is then

$$\mathbb{E}_{\psi}(p) = \sum_{k=-\infty}^{\infty} \hbar k |c_k|^2.$$

If an operator P satisfies $P\psi_k = \hbar k \psi_k$, then we could write $\mathbb{E}_{\psi}(p) = \langle P\psi, \psi \rangle$, the usual expression for the expectation value of an observable. One choice of such an operator is $P = -i\hbar \partial$, and we postulate that this operator is essentially the correct one on $L^2(\mathbb{R}^d)$ as well.

In the setting of Theorem 5.5, we obtain momentum operators by letting $q(x) = x_j$, yielding the operator¹² $P_j = -i \partial_j$. It is then self-adjoint with spectrum $\sigma(P_j) = \mathbb{R}$.

EXAMPLE 5.7: The Laplace operator and kinetic energy.

The kinetic energy T of a classical particle with mass m and velocity v is $T = mv^2/2 = p^2/(2m)$, where p = mv is its classical momentum. Substituting the momentum operator P_j for each coordinate of the classical momentum yields the operator $-\Delta/(2m)$.

We may obtain this from Theorem 5.5 via the polynomial $q(x) = ||x||^2$, yielding the kinetic energy operator¹³ – Δ . It is also self-adjoint with spectrum $\sigma(-\Delta) = [0, \infty)$. This is reasonable from a physical standpoint, in that kinetic energies cannot be negative but are presumably unbounded in principle.

¹² We choose units such that $\hbar = 1$.

¹³ Also choosing units such that m = 1/2.

The operator $-\Delta$ with domain \mathcal{D}_{max} is also called the *free Hamiltonian* ('free' since it describes the energy of a free particle) and is denoted H_0 .

PROPOSITION 5.8

If the polynomial q has real coefficients and g is a Borel function on $\sigma(q(D))$, then $g(q(D)) = \mathcal{F}^{-1}M_{g \circ q}\mathcal{F}$.

Furthermore, if $g \circ q \in L^2(\mathbb{R}^d)$, then g(q(D)) is a convolution operator on $L^2(\mathbb{R}^d)$ with kernel $\mathcal{F}^{-1}(g \circ q)$, i.e.,

$$g(q(D))f(x) = (2\pi)^{-d/2} \int_{\mathbb{R}^d} \mathcal{F}^{-1}(g \circ q)(x - y)f(y) \, \mathrm{d}y$$
 (5.2)

for all $f \in L^2(\mathbb{R}^d)$ and $x \in \mathbb{R}^d$.

PROOF. Let *E* be the spectral measure defined in the proof of Theorem 5.1, and let $E^q = q_*E$. Then

$$\int \mathrm{id}_{\mathbb{R}} \, \mathrm{d}E^q = \int q \, \mathrm{d}E = M_q,$$

and so E^q is the spectral resolution of M_q . It follows from Proposition 3.12 that $E^{q(D)} = \mathcal{F}^{-1}E^q\mathcal{F}$ is the spectral resolution of q(D), and so

$$\begin{split} g(q(D)) &= \int g \, \mathrm{d} E^{q(D)} = \mathcal{F}^{-1} \int g \, \mathrm{d} E^q \, \mathcal{F} \\ &= \mathcal{F}^{-1} \, \int g \circ q \, \mathrm{d} E \, \mathcal{F} = \mathcal{F}^{-1} M_{g \circ q} \mathcal{F}, \end{split}$$

proving the first claim.

For the second claim, notice that

$$g(q(D))f = \mathcal{F}^{-1}((g \circ q)\mathcal{F}f) = \mathcal{F}^{-1}(\mathcal{F}(\mathcal{F}^{-1}(g \circ q))\mathcal{F}f)$$
$$= (2\pi)^{-d/2}\mathcal{F}^{-1}(g \circ q) * f$$

by Proposition A.2, since if $g \circ q \in L^2(\mathbb{R}^d)$, then $\mathcal{F}^{-1}(g \circ q) \in L^2(\mathbb{R}^d)$. This proves (5.2).

EXAMPLE 5.9: The free propagator.

Let $\lambda \in \mathbb{C}$ with Re $\lambda > 0$ and consider the function $\varphi \colon \mathbb{R} \to \mathbb{C}$ given by $\varphi(x) = \mathrm{e}^{-\lambda x^2}$. Denote the Lebesgue measure on \mathbb{R}^d by m_d and write $m = m_1$. We aim to compute the Fourier transform of φ . First we evaluate the integral

$$I_1 = \int_{\mathbb{R}} \varphi(x) \, \mathrm{d}x.$$

Define $\eta: \mathbb{R}^2 \to \mathbb{R}$ by $\eta(x,y) = \sqrt{x^2 + y^2}$. It is well-known from Thorbjørnsen (2014, Opgave 11.4) that the image measure $\eta_* m_2$ has density $f(r) = 2\pi r \mathbf{1}_{(0,\infty)}(r)$ with respect to m. It follows from Fubini's theorem that

$$I_1^2 = \int_{\mathbb{R}^2} e^{-\lambda(x^2 + y^2)} dm_2(x, y) = \int_{\mathbb{R}^2} e^{-\lambda \eta(x, y)^2} dm_2(x, y)$$
$$= 2\pi \int_0^\infty r e^{-\lambda r^2} dr = \frac{\pi}{\lambda},$$

so $I_1 = \sqrt{\pi/\lambda}$.

Now let $c \in \mathbb{C}$ and consider the integral

$$I_2 = \int_{\mathbb{R}} \varphi(x+c) \, \mathrm{d}x. = \int_{\mathbb{R}} \mathrm{e}^{-\lambda(x+c)^2} \, \mathrm{d}x.$$

Writing c = a + ib, we let $R \ge |b|$ and consider the path $\Gamma(R)$ in $\mathbb C$ given by a rectangle with corners -R, R + ib and -R + ib that is traversed in the positive direction. Since φ extended to $\mathbb C$ is holomorphic, by Cauchy's integral theorem we get

$$0 = \oint_{\Gamma(R)} \varphi(z) dz = \int_{-R}^{R} \varphi(x) dx + \int_{0}^{b} \varphi(R + iy) dy$$
$$+ \int_{R}^{-R} \varphi(x + ib) dx + \int_{b}^{0} \varphi(-R + iy) dy.$$

As $R \to \infty$, the second and fourth integrals approach zero by the Dominated Convergence Theorem, so

$$\int_{\mathbb{R}} \varphi(x + i b) dx = \int_{\mathbb{R}} \varphi(x) dx = \sqrt{\frac{\pi}{\lambda}}.$$

Making the change of variables $x \mapsto x + a$ we obtain $I_2 = \sqrt{\pi/\lambda}$.

We are now ready compute the Fourier transform of φ :

$$\sqrt{2\pi}\widehat{\varphi}(t) = \int_{\mathbb{R}} e^{-\lambda x^2} e^{-itx} dx = e^{-t^2/4\lambda} \int_{\mathbb{R}} e^{-\lambda(x+it/2\lambda)^2} dx = \sqrt{\frac{\pi}{\lambda}} e^{-t^2/4\lambda}.$$

Now we pass to an arbitrary dimension d. First some notation: Let $q(x) = ||x||^2$ and $g(y) = e^{-\lambda y}$ for $x \in \mathbb{R}^d$ and $y \in \mathbb{R}$. Then $(g \circ q)(x) = e^{-\lambda ||x||^2}$. Writing $x = (x_1, \dots, x_d)$ and $t = (t_1, \dots, t_d)$ we get by Fubini's theorem

$$(2\pi)^{-d/2} \widehat{g \circ q}(t) = \int_{\mathbb{R}^d} e^{-\lambda ||x||^2} e^{-i t \cdot x} dx = \int_{\mathbb{R}^d} \prod_{j=1}^d e^{-\lambda x_j^2} e^{-i t_j x_j} dx$$
$$= (2\pi)^{-d/2} \prod_{j=1}^d \widehat{\varphi}(t_j) = \left(\frac{\pi}{\lambda}\right)^{d/2} \prod_{j=1}^d e^{-t_j^2/4\lambda}$$
$$= \left(\frac{\pi}{\lambda}\right)^{d/2} e^{-||t||^2/4\lambda},$$

implying that $\widehat{g \circ q}(t) = (2\lambda)^{-d/2} e^{-\|t\|^2/4\lambda}$. Since $g \circ q$ is an even function, $\mathcal{F}(g \circ q) = \mathcal{F}^{-1}(g \circ q)$. It follows from Proposition 5.8 that

$$e^{-\lambda H_0} \psi(x) = (4\pi\lambda)^{-d/2} \int_{\mathbb{R}^d} e^{-\|x-y\|^2/4\lambda} \psi(y) dy$$

for all $\psi \in L^2(\mathbb{R}^d)$.

We now find an expression for $\lambda = i$: First let $\psi \in L^1(\mathbb{R}^d) \cap L^2(\mathbb{R}^d)$. Then because $e^{-i(t-i\varepsilon)H_0} \psi \to e^{-itH_0} \psi$ in $L^2(\mathbb{R}^d)$ as $\varepsilon \to 0^+$ by the strong continuity of spectral integrals, there is a sequence (ε_n) approaching zero such that $e^{-i(t-i\varepsilon_n)H_0} \psi \to e^{-itH_0} \psi$ pointwise a.e. We then have the pointwise limit

$$\begin{split} \mathrm{e}^{-\mathrm{i}\,tH_0}\,\psi(x) &= \lim_{n \to \infty} \mathrm{e}^{-\mathrm{i}(t-\mathrm{i}\,\varepsilon_n)H_0}\,\psi(x) \\ &= \lim_{n \to \infty} (4\pi\,\mathrm{i}(t-\mathrm{i}\,\varepsilon_n))^{-d/2} \int_{\mathbb{R}^d} \mathrm{e}^{-\|x-y\|^2/4\mathrm{i}(t-\mathrm{i}\,\varepsilon_n)}\,\psi(y)\,\mathrm{d}y \\ &= (4\pi\,\mathrm{i}\,t)^{-d/2} \int_{\mathbb{R}^d} \mathrm{e}^{\mathrm{i}\|x-y\|^2/4t}\,\psi(y)\,\mathrm{d}y \end{split}$$

by the Dominated Convergence Theorem. Now let $\psi \in L^2(\mathbb{R}^d)$ and let $B_R \subseteq \mathbb{R}^d$ be the closed ball centred at 0 with radius $R \ge 0$. Then $\psi \mathbf{1}_{B_R} \in L^1(\mathbb{R}^d)$ and $\psi \mathbf{1}_{B_R} \to \psi$ in $L^2(\mathbb{R}^d)$ as $R \to \infty$. By the unitarity (and hence continuity) of $\mathrm{e}^{-\mathrm{i} t H_0}$ we then get

$$e^{-itH_0}\psi(x) = \lim_{R \to \infty} (4\pi i t)^{-d/2} \int_{B_R} e^{i||x-y||^2/4t} \psi(y) dy,$$

where the limit is taken in $L^2(\mathbb{R}^d)$. The kernel function

$$P(x,y;t) = (4\pi i t)^{-d/2} e^{i||x-y||^2/4t}$$
(5.3)

is called the *free propagator*. This has a straightforward physical interpretation: If ψ describes the initial state of a free particle, Proposition 4.5 implies that the state at time t is $e^{-itH_0}\psi$. Thus (5.3) says that the probability amplitude of finding the particle at x at time t is given by the amplitude of it starting at y, times the amplitude that it travelled from y to x, summed over all possible starting points. Then P(x,y;t) is the amplitude associated with starting at y and ending at x at time t. See Sakurai and Napolitano (2011, Section 2.6) for further discussion.

— 6 **—**

Conclusion and outlook

In this thesis we have covered some of the basic theory of unbounded operators, in particular self-adjoint operators, on Hilbert space, and discussed some of the applications to quantum mechanics. As any student of quantum mechanics well knows, this is just the beginning, and there are many directions we could go in from here.

One possibility is to look at some of the topics usually covered in an undergraduate course in quantum mechanics: The harmonic oscillator, the hydrogen atom, the uncertainty principle, spin, etc. For a rigorous reference on these topics as well as more advanced topics let us mention in particular Hall (2013).

Another is to look closer at the mathematical framework. In Subsection 5.2 we showed that the free Hamiltonian H_0 is self-adjoint, and one might for instance wonder how to deal with more general Schrödinger operators in which the potential is nonzero. One answer is given by the *Kato-Rellich the-orem*, which ensures the self-adjointness of certain 'small' perturbations of self-adjoint operators. A *Kato-Rellich potential* is a Borel function $V: \mathbb{R}^d \to \mathbb{C}$ such that either $V \in L^2(\mathbb{R}^d) + L^\infty(\mathbb{R}^d)$ for d = 1, 2, 3, or $V \in L^p(\mathbb{R}^d) + L^\infty(\mathbb{R}^d)$ for some p > d/2 if $d \ge 4$. We then have the following corollary to the Kato-Rellich theorem:

COROLLARY 6.1

If V is a real-valued Kato–Rellich potential and M_V is the corresponding multiplication operator, then $H_0 + M_V$ is self-adjoint on $\mathcal{D}(H_0)$.

We refer to Schmüdgen (2012, Chapter 8) for the precise statement and proof of the Kato–Rellich theorem, the proof of the above corollary, as well as further topics on perturbations of self-adjoint operators.

— A —

The Fourier transform and tempered distributions

We use the symbol ∂_j to denote differentiation with respect to the jth variable. For derivatives of higher-order we use multi-index notation. A multi-index is a d-tuple $\alpha \in \mathbb{N}_0^d$ of non-negative integers. For a multi-index $\alpha = (\alpha_1, \ldots, \alpha_d)$ and a point $x = (x_1, \ldots, x_d) \in \mathbb{R}^d$ we let

$$|\alpha| = \sum_{j=1}^d \alpha_j$$
, $\partial^{\alpha} = \partial_1^{\alpha_1} \cdots \partial_d^{\alpha_d}$, and $x^{\alpha} = \prod_{j=1}^d x_j^{\alpha_j}$.

We will usually apply the operator ∂^{α} to smooth functions, so the order of differentiation is immaterial. It is also convenient to introduce the differential operator $D^{\alpha} = (-\mathrm{i})^{|\alpha|} \partial^{\alpha}$. For convenience, we use the symbols x^{α} and t^{α} to also denote the functions $x \mapsto x^{\alpha}$ and $t \mapsto t^{\alpha}$.

A.1 • Tempered distributions

We first introduce the *Schwartz space* S, also called the space of *rapidly decreasing functions*: This is the subspace of $C^{\infty}(\mathbb{R}^d)$ of functions φ such that $x^{\alpha}D^{\beta}\varphi$ is bounded for all multi-indices α and β . The Schwartz space becomes a Fréchet space when equipped with an appropriate family of seminorms, though we shall not be concerned with the details. Notice that $S \subseteq L^p(\mathbb{R}^d)$ for all $p \in [1, \infty]$.

The elements of the topological dual space S' are called *tempered distributions* and S' is equipped with the weak*-topology.

If $F \in \mathcal{S}'$ and $\varphi \in \mathcal{S}$, then we write $\langle F, \varphi \rangle = F(\varphi)$. If further there is some $f \in L^p(\mathbb{R}^d)$ for $p \in [1, \infty]$ such that $\langle F, \varphi \rangle = \int f(x)\varphi(x)\,\mathrm{d}x$, then we say that F is given by f. Notice that the integral converges by the Hölder inequality. It turns out that the map $f \mapsto F$ is injective (see for instance Grubb (2009, Lemma 3.2)), so we may identify F and f.

The subspace $\mathcal{O}_M(\mathbb{R}^d) \subseteq C^\infty(\mathbb{R}^d)$ of *slowly increasing functions* consists of those $\psi \in C^\infty(\mathbb{R}^d)$ that satisfy the following property: For all $\alpha \in \mathbb{N}_0^d$ there

exist C > 0 and $N \in \mathbb{N}$ such that

$$|\partial^{\alpha}\psi(x)| \le C(1+||x||)^N$$
, $(x \in \mathbb{R}^d)$.

In particular, any polynomial is slowly increasing. We also see that multiplication by elements in $\mathcal{O}_M(\mathbb{R}^d)$ preserves \mathcal{S} . Given $F \in \mathcal{S}'$ and $\psi \in \mathcal{O}_M(\mathbb{R}^d)$, we define ψF by $\langle \psi F, \varphi \rangle = \langle F, \psi \varphi \rangle$ for all $\varphi \in \mathcal{S}$. We also define the derivative $\partial^\alpha F$ by $\langle \partial^\alpha F, \varphi \rangle = (-1)^{|\alpha|} \langle F, \partial^\alpha \varphi \rangle$, and similarly for the operator D^α . In each case we again obtain tempered distributions.

We refer to Grubb (2009, Chapter 5) for details and further results.

A.2 • The Fourier transform

We assume familiarity with the Fourier transform on $L^1(\mathbb{R}^d)$ and $L^2(\mathbb{R}^d)$ for d = 1. Many of the properties of the Fourier transform generalise to arbitrary $d \in \mathbb{N}$. For example, L^1

$$(it)^{\alpha} \partial^{\beta} \mathcal{F} f = \mathcal{F} (\partial^{\alpha} (-ix)^{\beta} f), \tag{A.1}$$

where $\mathcal{F}f = \hat{f}$ is the Fourier transform of $f \in L^1(\mathbb{R}^d)$, cf. Folland (2007, Theorem 8.22). Instead using the operator D, this formula takes the simpler form

$$t^{\alpha}D^{\beta}\mathcal{F}f = \mathcal{F}(D^{\alpha}(-x)^{\beta}f). \tag{A.2}$$

THEOREM A.1

The d-dimensional Fourier transform is an isomorphism of S onto itself.

Thus we may extend the Fourier transform to a linear map on S' by defining $\langle \hat{F}, \varphi \rangle = \langle F, \hat{\varphi} \rangle$ for $F \in S'$ and $\varphi \in S$. The basic properties of the Fourier transform extend directly to this setting. For instance, the identity (A.2) becomes

$$t^{\alpha}D^{\beta}\mathcal{F}F = \mathcal{F}(D^{\alpha}(-x)^{\beta}F),\tag{A.3}$$

cf. Folland (2007, p. 295).

The convolution theorem also generalises to arbitrary d in the following way:

¹⁴ To avoid writing too many parentheses, we adopt the convention that expressions such as $\partial^{\alpha} x^{\beta} f$ are right-associative.

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PROPOSITION A.2

For $f, g \in L^2(\mathbb{R}^d)$, we have

$$\mathcal{F}^{-1}(\mathcal{F}f\cdot\mathcal{F}g)=(2\pi)^{-d/2}f*g.$$

PROOF. Lieb and Loss (2001, Theorem 5.8).

— B —

Miscellaneous functional analysis

PROPOSITION B.1

The space $C_c^{\infty}(\mathbb{R}^d)$ (and hence also S) is dense in $L^p(\mathbb{R}^d)$ for $p \in [1, \infty)$.

PROOF. Folland (2007, Proposition 8.17).

THEOREM B.2: The Bounded Linear Transformation Theorem

Let X be a normed space and Y a Banach space, and let $T: U \to Y$ be a bounded linear map defined on a subspace U of X. Then T can be uniquely extended to a bounded operator $\tilde{T}: \overline{U} \to Y$ on the closure of U. Furthermore, $||\tilde{T}|| = ||T||$.

PROOF. Let $x \in \overline{U}$, and let $(x_n)_{n \in \mathbb{N}}$ be a sequence in U that converges to x. Then

$$||Tx_n - Tx_m|| \le ||T|| \cdot ||x_n - x_m||,$$

and the right-hand side converges to zero, so $(Tx_n)_{n\in\mathbb{N}}$ is a Cauchy sequence. By the completeness of Y, this sequence converges to some $y\in Y$. Now let $(x_n')_{n\in\mathbb{N}}$ be another sequence that converges to x, and let $y'=\lim_{n\to\infty}Tx_n'$. Then

$$||y - y'|| \le ||T|| \cdot ||x_n - x_n'||$$

and the right-hand side converges to zero, so y = y'. Thus we may define $\tilde{T}x = y$.

We now show that \tilde{T} is bounded. Let $x \in \overline{U}$ and let the sequence $(x_n)_{n \in \mathbb{N}}$ converge to x. Notice that

$$\|\tilde{T}x\| = \lim_{n \to \infty} \|Tx_n\| \le \|T\| \lim_{n \to \infty} \|x_n\| = \|T\| \cdot \|x\|,$$

showing that $\|\tilde{T}\| \le \|T\|$, so \tilde{T} is bounded. The opposite inequality is obvious, so in total $\|\tilde{T}\| = \|T\|$.

B.1 • Direct sums of vector spaces

The *direct sum* $X \oplus Y$ of vector spaces X and Y consists of the set $X \times Y$ equipped with pointwise addition and scalar multiplication. If X and Y are normed spaces, we equip $X \oplus Y$ with a norm $\|\cdot\|$ given by $\|(x,y)\|^2 = \|x\|_X^2 + \|y\|_Y^2$. If X and Y are Banach spaces, then it is easy to see that $X \oplus Y$ is also a Banach space if equipped with the above norm.

If $\mathcal H$ and $\mathcal K$ are Hilbert spaces, then we equip their direct product $\mathcal H\oplus\mathcal K$ as vector spaces with an inner product

$$\langle (x_1, y_1), (x_2, y_2) \rangle = \langle x_1, x_2 \rangle_{\mathcal{H}} + \langle y_1, y_2 \rangle_{\mathcal{K}}.$$

Note that the induced norm on $\mathcal{H} \oplus \mathcal{K}$ agrees with the above norm inherited from the norms on \mathcal{H} and \mathcal{K} . Thus $\mathcal{H} \oplus \mathcal{K}$ is in fact a Hilbert space.

It is possible to define the direct sum of an arbitrary collection of Hilbert spaces, though we shall only need finite sums. See e.g. Conway (1990, Section I.6) for the case of a sequence of Hilbert spaces.

B.2 • Bounded sesquilinear forms

Our proof of Stone's theorem (cf. Theorem 4.4) uses a fundamental property of bounded sesquilinear forms which we now recall. A bounded sesquilinear form on a Hilbert space $\mathcal H$ is a sesquilinear map $\Gamma\colon\mathcal H\times\mathcal H\to\mathbb C$ such that the norm

$$||\Gamma|| = \sup \{\Gamma(x, y) \mid x, y \in \mathcal{H}, ||x|| = ||y|| = 1\}$$

is finite. The following result relates bounded sesquilinear forms and bounded operators on \mathcal{H} :

THEOREM B.3

Let $\Gamma: \mathcal{H} \times \mathcal{H} \to \mathbb{C}$ be a bounded sesquilinear form. Then there exists an unique $T \in \mathcal{B}(\mathcal{H})$ such that

$$\Gamma(x,y) = \langle Tx, y \rangle, \quad (x,y \in \mathcal{H}).$$

PROOF. Rudin (1991, Theorem 12.8), Conway (1990, Theorem II.2.2) or Skibsted (2019, Theorem 3.20).

B.3 • Consequences of the Baire category theorem

In this section we recall some important corollaries of the Baire category theorem that we will need. We assume familiarity with the category theorem and refer to Folland (2007, Section 5.3) for details. The theorems below generalise significantly, though we shall not need these generalisations. Discussion of this can be found e.g. in Rudin (1991, Chapter 2).

THEOREM B.4: The Bounded Inverse Theorem

If $T: X \to Y$ is a continuous linear bijection between Banach spaces, then T^{-1} is continuous.

THEOREM B.5: The Closed Graph Theorem

Let X and Y be Banach spaces, and let $T: X \to Y$ be a closed linear map. Then T is bounded.

− ^c − Spectral theory

Throughout this section, (Ω, \mathcal{E}) denotes a measurable space and \mathcal{H} denotes a non-zero complex Hilbert space.

C.1 • Spectral measures

Our definition of spectral measures is as follows:

DEFINITION C.1: Spectral measures

A *spectral measure* on \mathcal{E} is a map $E \colon \mathcal{E} \to \mathcal{B}(\mathcal{H})$ which satisfies the following properties:

- (i) $E(\Omega) = I$.
- (ii) E(A) is an orthogonal projection for all $A \in \mathcal{E}$.
- (iii) *E* is countably additive, i.e. for any sequence $(A_n) \subseteq \mathcal{E}$ of pairwise disjoint

sets,

$$E\left(\bigcup_{n\in\mathbb{N}}A_n\right) = \operatorname{s-lim}_{N\to\infty}\sum_{n=1}^N E(A_n).$$

The presentation of spectral measures in the present thesis is based on Schmüdgen (2012, Section 4.2). We remark that our definition of spectral measures in Definition C.1 is the same as the definition of *projection-valued measures* in Reed and Simon (1980, Section VII.3). Some authors, e.g. Rudin (1991, Definition 12.17), also use the term *resolution of the identity*, though note that Schmüdgen (2012, Definition 4.1) uses this term differently¹⁵.

It follows by a standard argument that spectral measures are strongly continuous from above and below. They are also clearly finitely additive.

Another important property of spectral measures is that we can characterise them using positive measures: It follows readily from the the definition that the map $E_x \colon \mathcal{E} \to [0, \infty)$ given by $E_x(A) = \langle E(A)x, x \rangle$ is a (positive) measure on \mathcal{E} for all $x \in \mathcal{H}$. This property is actually equivalent to property (iii), assuming the other two properties. In fact, many authors take this characterisation as part of the definition of spectral measures in place of (iii). This equivalence is the content of Schmüdgen (2012, Lemma 4.4). Our definition is thus also equivalent to the definition of spectral measures in Skibsted (2003, Definition 1.1), and the definition of resolutions of the identity in Skibsted (2019, Section 3.4).

Rudin (1991, Definition 12.17) substitutes the assumption of countable additivity for the property that the maps $E_{x,y} \colon \mathcal{E} \to \mathbb{C}$ given by $E_{x,y}(A) = \langle E(A)x,y \rangle$ is a complex measure¹⁶ for all $x,y \in \mathcal{H}$. But this is clearly equivalent to the above property concerning positive measures. Rudin also requires that a spectral measure be finitely additive and multiplicative in the following sense:

PROPOSITION C.2: Multiplicativity of spectral measures

Let E be a spectral measure on E. For all $A, B \in \mathcal{E}$ we have

$$E(A)E(B) = E(A \cap B).$$

PROOF. Schmüdgen (2012, Lemma 4.3) or Skibsted (2019, Remark p. 82). □

¹⁵ To wit, a resolution of the identity for Schmüdgen is a one-parameter family $E: \mathbb{R} \to \mathcal{B}(\mathcal{H})$ of orthogonal projections that is monotone, strongly right continuous and with the property that $\lim_{\lambda \to -\infty} E(\lambda)x = 0$ and $\lim_{\lambda \to +\infty} E(\lambda)x = x$ for $x \in \mathcal{H}$. There is, however, a one-to-one correspondence between resolutions of the identity and spectral measures on $\mathfrak{B}(\mathbb{R})$, see Schmüdgen (2012, Theorem 4.6).

¹⁶ For the theory of complex measures we refer to Rudin (1986) and Folland (2007).

Thus our definition of spectral measures agrees with Rudin's definition of resolutions of the identity. However, in the case where X is a compact or locally compact Hausdorff space and $\mathcal{E} = \mathfrak{B}(X)$, Rudin assumes that the measures $E_{x,y}$ are *regular* Borel measures. We have no need for this assumption, cf. Skibsted (2019, Remark 3.61).

It is worth remarking on the multiplicativity of spectral integrals. A positive (or complex) measure μ whose corresponding integral is multiplicative would have to satisfy $\mu(A)^2 = \mu(A)$ for all $A \in \mathcal{E}$, since $\mathbf{1}_A \cdot \mathbf{1}_A = \mathbf{1}_A$. But then either $\mu(A) = 0$ or $\mu(A) = 1$, which is not very interesting. However, there are many more idempotents in $\mathcal{B}(\mathcal{H})$, allowing for a useful multiplicative integral.

We also note that we cannot (in general) hope for convergence of the sum $\sum_{n=1}^{\infty} E(A_n)$ in Definition C.1(iii) with respect to the operator norm. If (P_n) is any sequence of (not necessarily orthogonal) projections, then the sum $\sum_{n=1}^{\infty} P_n$ converges in the operator norm if and only if finitely many P_n are nonzero. For if P is idempotent and $Px \neq 0$ for some $x \in \mathcal{H}$, then $||Px|| = ||P^2x|| \leq ||P|| \cdot ||Px||$, so $||P|| \geq 1$.

C.2 • The spectral theorem for a bounded normal operator

THEOREM C.3: Existence of spectral integrals of $L^{\infty}(E)$ -functions

There exists an isometric *-homomorphism $\Psi: L^{\infty}(E) \to \mathcal{B}(\mathcal{H})$ *such that*

$$\langle \Psi(f)x, x \rangle = \int_{\Omega} f \, \mathrm{d}E_x$$

for all $f \in L^{\infty}(E)$ and $x \in \mathcal{H}$. This justifies the notation $\Psi(f) = \int_{\Omega} f \, dE$.

PROOF. Rudin (1991, Theorem 12.21) or Skibsted (2019, Theorem 3.56). □

THEOREM C.4: The spectral theorem for a bounded normal operator

Let T be a bounded normal operator on H. There exists a unique resolution E of the identity on the Borel σ -algebra $\mathfrak{B}(\sigma(T))$ on the spectrum of T such that

$$T = \int_{\sigma(T)} \lambda \, \mathrm{d}E(\lambda).$$

The resolution E of the identity is called the spectral resolution *of T*.

PROOF. Rudin (1991, Theorem 12.23) or Skibsted (2019, Theorem 3.65).

List of symbols

In the list below, X denotes a topological space, \mathcal{E} is a σ -algebra on some set, Ω is a subset of \mathbb{R}^d (which is often open), V and W are vector space (equipped with appropriate structure), T is a linear map, \mathcal{H} is a Hilbert space, and E is a spectral measure.

Note that we dot attempt do distinguish between spaces of functions and spaces of equivalence classes of functions, e.g. between $\mathcal{L}^p(\mu)$ and $L^p(\mu)$.

Measure theory

- $\mathfrak{B}(X)$ The Borel σ -algebra on X.
- $\int_A f(t) dt$ The integral of a Lebesgue-measurable function f on the domain $A \subseteq \mathbb{R}$.
- $\mathbf{1}_A$ The indicator/characteristic function on a set A.

Function spaces

- $\mathcal{M}(\mathcal{E})$ The complex-valued $(\mathcal{E}, \mathfrak{B}(\mathbb{C}))$ -measurable functions.
- $\mathcal{SM}(\mathcal{E})$ The subalgebra of $\mathcal{M}(\mathcal{E})$ of simple functions.
- C(X) The continuous functions on X.
- $C_c(X)$ The continuous functions on X with compact support.
- $C^k(\Omega)$ The *k* times continuously differentiable functions on Ω .
- $C_0^k(\mathbb{R})$ The subspace of $C^k(\mathbb{R})$ of functions that go to zero at $\pm \infty$.
- $C_c^{\infty}(\mathbb{R})$ The smooth functions on \mathbb{R} with compact support.
- $L^{\infty}(E)$ The space of functions that are essentially bounded with respect to E.
- $L^p(\Omega)$ The space of *p*-integrable functions with respect to the Lebesgue measure on Ω .
- S, S' The Schwartz space and its topological dual.
- $\mathcal{O}_M(\mathbb{R}^d)$ The space of slowly increasing functions on \mathbb{R}^d .

Linear algebra

- $\mathcal{B}(V, W)$ The vector space of everywhere-defined bounded linear operators from V to W. If V = W we also write $\mathcal{B}(V, V) = \mathcal{B}(V)$.
- $\mathcal{R}(T)$ The range of T.

continues on the next page

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- $\mathcal{N}(T)$ The kernel of T.
- G(T) The graph of T.
- $\langle \cdot, \cdot \rangle_V$ The inner product on V. We omit the subscript if V is clear from context. Our inner products are linear in the first argument.
- $U(\mathcal{H})$ The group of unitary operators on \mathcal{H} .

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