

Functional Analysis

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Part I

Vector Spaces with a Metric

Functional analysis is the interlace of algebra and analysis, in which algebraic structures are endowed with topological structure. The approach's utility counts for the rapid growth of applications over the past century, be it in quantum mechanics, statistics, or computing science. Rarely are we concerned with a single object, like a function, a random variable, or a measure, but instead consider large classes of such objects. One technique for handling these classes is to add algebraic structure which elaborates on the natural relation between these objects. Functional analysis adds in analytical structure to the class, so we may combine the algebraic and analytical structure to obtain powerful results about the objects involved.

Example. *We rarely analyze a measurable function f in isolation. Instead, we prove theorems about a class of measurable functions defined on the same measurable space. If f and g are measurable, then we may consider their addition $f + g$, their multiplication fg , and scaling λf (for $\lambda \in \mathbf{R}$), which are all measurable. Thus the space of measurable functions on a set is a vector space. Similarly, the sum and product of two continuous functions on a topological space is continuous, so, we may consider $C[0, 1]$, the space of all continuous, real-valued functions on the unit interval $[0, 1]$, as a vector space.*

The most common analytic operations are scaling and rotation. As a result, the most common vector spaces which occur in functional analysis are defined over the real and complex numbers. More specifically, real vector spaces X are those equipped with natural scaling operations $x \mapsto \lambda x$ for $\lambda > 0$, combined with ‘reflection’ operators $x \mapsto -x$. If X is a real vector space with a specified ‘counterclockwise twisting’, i.e. a linear map $J : X \rightarrow X$ with $J^2(x) = -x$ for all $x \in X$, then there is a natural way to turn X into a complex vector space by defining $(a + ib)x = ax + b(Jx)$. The geometry of the complex numbers then enables us to define a counterclockwise twist by any angle θ , by the operation $x \mapsto e^{i\theta}x$. Thus we have a natural operation of the complex numbers over the space.

Chapter 1

Basics of Banach Spaces

The most natural way to associate a space with topological structure is by associating a metric on the space. There is an even more natural vector space operation on a vector space. We can think of elements of a vector space as ‘arrows’ extending from the origin. Thus a natural way to add a topology is to give these arrows a ‘length’, and the function associating a vector with a length is known as a *norm*. More precisely, a *pre-norm* on a vector space X is a map $\|\cdot\| : X \rightarrow [0, \infty)$ which is *homogenous*, in the sense that for each $\alpha \in \mathbf{K}$ and $x \in X$, $\|\alpha x\| = |\alpha| \|x\|$, and satisfies the *triangle inequality* $\|x + y\| \leq \|x\| + \|y\|$ for each $x, y \in X$. If in addition, we have $\|x\| = 0$ if and only if $x = 0$, then we say $\|\cdot\|$ is a *norm*. A vector space X equipped with a norm is called a *norm space*.

On finite dimensional spaces, there is essentially a unique way to define a norm on the space. We say two norms $\|\cdot\|_1$ and $\|\cdot\|_2$ on a vector space X are *comparable* if for each $x \in X$, $\|x\|_1 \sim \|x\|_2$. It is a simple consequence of the Heine-Borel theorem that all norms on a finite dimensional space are normable. Thus there is only a single notion of ‘size’ on such spaces. On the other hand, there are many ways to measure the size of a function, or more general vectors in infinite dimensional norm spaces.

Example. Let X be a measure space. For each $0 < p < \infty$, let $\mathcal{L}^p(X)$ denote the space of all measurable functions $f : X \rightarrow \mathbf{K}$ such that

$$\|f\|_p = \left(\int |f(x)|^p dx \right)^{1/p} < \infty.$$

Minkowski’s inequality tells us that $\|\cdot\|_p$ is a prenorm on $\mathcal{L}^p(X)$. Notice that if u is a function such that $u(x) = 0$ for almost every x , then $\|f\|_p = \|f + u\|_p$.

In particular, this means that $\|\cdot\|_p$ is a well defined operation on the quotient space of $\mathcal{L}^p(X)$ by the vector space of functions which are equal to zero almost everywhere. We denote this quotient space by $L^p(X)$. Since $\|f\|_p = 0$ if and only if f is equal to zero almost everywhere, this implies that $L^p(X)$, equipped with the induced norm $\|\cdot\|$, is a norm space.

Example. Given a measurable function f on a measure space X , we can define the L^∞ pre-norm

$$\|f\|_\infty = \inf\{t > 0 : |f(x)| \leq t \text{ for almost every } x\}.$$

The space of functions with finite L^∞ pre-norm is denoted by $\mathcal{L}^\infty(X)$, and the quotient space by the class of functions equal to zero is denoted $L^\infty(X)$.

Example. On a compact topological space K , we can consider the family of all bounded continuous functions from K to \mathbf{K} , denoted $C(K)$. The quantity

$$\|f\|_\infty = \sup\{|f(x)| : x \in K\},$$

gives $C(K)$ the structure of a norm space. More generally, if X is any norm space, we can consider the family of all bounded continuous functions from K to X , with associated norm

$$\|f\|_{C(K,X)} = \sup\{\|f(x)\| : x \in K\}.$$

If X is a Banach space, then $C(K, X)$ is a Banach space.

For the purpose of the general theory, one need not know too much measure theory. The most important norm spaces for intuition are the simplest infinite dimensional examples of the spaces $L^p(X)$, where X is the collection of all natural numbers, with the associated counting measure. This space is often denoted l^p , and can be practically defined as the collection of all sequences $\{a_n\}$ such that

$$\|a\|_p = \left(\sum_{n=1}^{\infty} |a_n|^p \right)^{1/p} < \infty.$$

For $p = \infty$, we let l_∞ denote the collection of all bounded sequences, and let $\|a\|_\infty = \sup |a_n|$.

Remark. We will focus on norms in these notes, but prenorms are really no more general. If X is a space with a prenorm $\|\cdot\|$, then the space $X_0 = \{x \in X : \|x\| = 0\}$ forms a subspace of X . We have $\|x\| = \|x + x_0\|$ for each $x_0 \in X_0$, so if we form the quotient space $X' = X/X_0$, then the quantity $\|x + X_0\|$ is well defined, and this is a *norm* on X' . So results about norms automatically extend to results about prenorms where we identify vectors which differ by an element of X_0 in a single equivalence class. This is precisely what we did to obtain the norm spaces $L^p(X)$ from the prenorm spaces $\mathcal{L}^p(X)$.

For each norm on a space X , we can associate a metric given by $d(x, y) = \|x - y\|$. The most powerful theorems of functional analysis deal with the study of norm spaces where the induced metric on a norm space is *complete*. We say such a space is a *Banach space*. It is a basic fact of measure theory that the spaces $L^p(X)$ are complete, and it is also easy to see that the space $C(K)$ is complete, as is proved in a basic class in real analysis.

Theorem 1.1. $L^p(X)$ is a Banach space.

Proof. We leave the $p = \infty$ case to the reader, since the argument is elementary. Let f_1, f_2, \dots be a Cauchy sequence in $L^p(X)$. Then we can find a subsequence f_{k_i} such that for each $m \geq 1$,

$$\|f_{k_{i+m}} - f_{k_i}\|_p \leq 1/4^k.$$

In particular, this means that the infinite sum

$$\sum_{i=1}^{\infty} |f_{k_{i+1}}(x) - f_{k_i}(x)|$$

converges absolutely for almost all x . This is a consequence of the Borel-Cantelli lemma, since if

$$E_k = \{x : |f_{k_{i+1}}(x) - f_{k_i}(x)| \geq 1/2^{pk}\},$$

then Markov's inequality implies $|E_i| \leq 1/2^{ip}$, and so $\sum_{i=1}^{\infty} |E_i| < \infty$. Thus we can define a measurable function f almost everywhere by the infinite series

$$f(x) = f_{k_1}(x) + \sum_{i=1}^{\infty} f_{k_{i+1}}(x) - f_{k_i}(x) = \lim_{n \rightarrow \infty} f_{k_n}(x).$$

Fatou's lemma and Minkowski's inequality imply that

$$\begin{aligned}\|f - f_{k_i}\|_p &= \left(\int \left| \lim_{n \rightarrow \infty} f_{k_n}(x) - f_{k_i}(x) \right|^p dx \right)^{1/p} \\ &\leq \liminf_{n \rightarrow \infty} \left(\int |f_{k_n}(x) - f_{k_i}(x)|^p dx \right)^{1/p} \\ &\leq \liminf_{n \rightarrow \infty} \|f_{k_n} - f_{k_i}\|_p \leq 1/4^k.\end{aligned}$$

Thus $f_{k_i} \rightarrow f$ in the L^p norm. Since $\{f_k\}$ is a Cauchy sequence, it is easy to justify that it is also true that $f_k \rightarrow f$ in the L^p norm. Thus Cauchy sequences converge in $L^p(X)$. \square

Remark. Studying Banach spaces is not too much more general than studying norm spaces. Given any norm space X , we can always find a unique Banach space X' containing X , such that X is dense in X' . The space X' is known as the *completion* of X . One can view the elements of $X' - X$ as 'artificial' elements of X . But, nonetheless are key to understanding the geometric structure of X . For instance, if we equip $C[0, 1]$ with the L^1 norm induced by the Lebesgue measure, then $C[0, 1]$ is an incomplete norm space. The completion of $C[0, 1]$ with respect to this norm is the space $L^1[0, 1]$ of integrable functions. One perspective of the class of measurable functions is that they are an artifice which provides us with stronger tools to analyze the class of the more basic functions which we grew familiar with in our first years of university.

The examples above provide most of the intuition we will need for the basic Banach space theory. But it is useful to know other examples for applications in other areas of analysis.

Example. Let K be a compact metric space. For each $0 < p \leq 1$, we say a function $f : K \rightarrow \mathbf{K}$ is Lipschitz of order p if $|f(x) - f(y)| \lesssim d(x, y)^p$ for each $x, y \in K$. The set of such functions forms a vector space $\text{Lip}_p(K)$, which is a subspace of $C[0, 1]$, but not a closed subspace. For each $f \in \text{Lip}_p[0, 1]$, we define the Lipschitz norm

$$\|f\|_{\text{Lip}_p(K)} = \|f\|_\infty + \sup \frac{|f(x) - f(y)|}{d(x, y)^p}.$$

This norm gives $\text{Lip}_p(K)$ the structure of a complete metric space. If $\{f_n\}$ is a Cauchy sequence in $\text{Lip}_p(K)$, the completeness of $L^\infty(K)$ shows there exists

$f \in L^\infty(K)$ such that $f_n \rightarrow f$ in the L^∞ norm. Since $\{f_n\}$ is Cauchy, there exists a constant M such that for all n , $|f_n(x) - f_n(y)| \leq M|x - y|^p$. In particular, the uniform convergence of f_n to f implies that $|f(x) - f(y)| \leq M|x - y|^p$, so $f \in \text{Lip}_p(K)$. Now

$$\begin{aligned} & \lim_{n \rightarrow \infty} \sup_{x, y \in K} \frac{|[f(x) - f_n(x)] - [f(y) - f_n(y)]|}{d(x, y)^p} \\ &= \lim_{n \rightarrow \infty} \sup_{x, y \in K} \lim_{m \rightarrow \infty} \frac{[f_m(x) - f_n(x)] - [f_m(y) - f_n(y)]}{d(x, y)^p} \\ &\leq \lim_{n \rightarrow \infty} \limsup_{m \rightarrow \infty} \sup_{x, y \in K} \frac{[f_m(x) - f_n(x)] - [f_m(y) - f_n(y)]}{d(x, y)^p} \\ &\leq \lim_{n \rightarrow \infty} \limsup_{m \rightarrow \infty} \|f_m - f_n\|_{\text{Lip}_p(K)} = 0. \end{aligned}$$

The last equality is equivalent to the fact that $\{f_n\}$ is a Cauchy sequence in $\text{Lip}_p(K)$. But this means that f_n converges to f in the $\text{Lip}_p(K)$ metric. Thus $\text{Lip}_p(K)$ is a complete metric space.

We can develop the theory of series in a norm space. Given a sequence $\{x_n\}$. We say the sequence is absolutely summable if $\sum \|x_n\| < \infty$. The same proof as for \mathbf{R} shows that such a sequence is *unconditionally summable*, i.e. that for each permutation $\pi : \mathbf{N} \rightarrow \mathbf{N}$,

$$\sum_{n=1}^{\infty} x_{\pi(n)} = \sum_{n=1}^{\infty} x_n.$$

A space X is a Banach space if and only if every absolutely summable series converges, which was a fact we implicitly used in our proof that $L^p(X)$ was a Banach space. However, in infinite dimensional spaces, there are always examples of unconditionally summable sequences which are not absolutely summable. This is a result of A. Dvoretzky and C.A. Rogers.

Example. Consider the space c_0 . For each n , let e_n denote the sequence which is equal to zero except at the n 'th position, which takes the value 1. If we set $x_n = e_n/n$, then the sequence $\{x_n\}$ is not absolutely summable, since

$$\sum_{n=1}^{\infty} \|x_n\|_{\infty} = \sum_{n=1}^{\infty} \frac{1}{n} = \infty.$$

Nonetheless, the sequence is unconditionally summable to the sequence $\sum_{n=1}^{\infty} e_n/n$.

1.1 Subspaces

If Y is a subspace of a norm space X , then Y obviously has the natural structure of a norm space. We note that if X is a Banach space, then Y is a Banach space under this induced norm if and only if Y is a *closed* subspace of X . Here are some examples.

Example. Let $c_0 \subset l_\infty$ be the family of all sequences $\{x_n\}$ such that $\lim_{n \rightarrow \infty} x_n = 0$. Then c_0 is a closed subspace of l_∞ , and is thus a Banach space. The space c_{00} of all sequences $\{x_n\}$ which are only non-zero for finitely many values is a subspace of c_0 , but it is not closed, and therefore not a Banach space. c_{00} is dense in l_p for all $1 \leq p < \infty$, and dense in c_0 , but is not dense in l_∞ .

Example. Let \mathbf{T} be the interval $[-\pi, \pi)$ equipped with the Lebesgue measure. For each $f \in L^p(\mathbf{T})$, and $n \in \mathbf{Z}$, we define

$$\hat{f}(n) = \frac{1}{2\pi} \int_{-\pi}^{\pi} f(x) e^{-nix} dx.$$

Hölder's inequality shows that $|\hat{f}(n)| \leq \|f\|_p$ for each $n \in \mathbf{Z}$. In particular, this means the set

$$H^p(\mathbf{T}) = \{f \in L^p(\mathbf{T}) : \hat{f}(n) = 0 \text{ if } n < 0\}$$

is closed, and is therefore a Banach space. These are examples of Hardy spaces on \mathbf{T} . They can be identified with holomorphic functions on the unit disk under the identification of f with the analytic function

$$u(z) = \sum_{n=0}^{\infty} \hat{f}(n) z^n,$$

which converges for $|z| < 1$ since $\{\hat{f}(n)\}$ is a bounded sequence. The L^p norm also implies u satisfies a suitable growth condition as z approaches the boundary of the unit disk, but we do not dwell on this fact.

Before we move on, let us introduce some notation. For any norm space X , we let B_X denote the closed unit ball in X , and S_X the unit sphere, i.e.

$$B_X = \{x \in X : \|x\| \leq 1\} \quad \text{and} \quad S_X = \{x \in X : \|x\| = 1\}.$$

Given two sets A and B , we let $A + B = \{a + b : a \in A, b \in B\}$. Then for each $a \in X$, $a + B_X$ is the closed unit ball centered at a . Similarly, we can define $A - B = \{a - b : a \in A, b \in B\}$, and for each $\alpha \in \mathbf{K}$, $\alpha A = \{\alpha a : a \in A\}$.

1.2 Convexity

A critical notion in functional analysis is *convexity*. A set $E \subset X$ is *convex* if for each $x, y \in E$ and $\lambda \in (0, 1)$, $\lambda x + (1 - \lambda)y \in E$. A set E is *balanced* if $\alpha E \subset E$ for each $|\alpha| \leq 1$, and *symmetric* if $-E = E$. And a set E is *absorbing* if, for each $x \in X$, there exists t such that for $|\alpha| \geq t$, $x \in \alpha E$. As an example, the unit ball B_X is closed, convex, and absorbing. In the future, we will make heavy use of the following result.

Theorem 1.2. *Every closed, convex, absorbing set in a Banach space contains a neighbourhood of the origin.*

Proof. Let $E \subset X$ be a closed, convex, absorbing set. Then $E \cap (-E)$ is closed, convex, and absorbing, so we may assume without loss of generality that E is symmetric. It suffices to show E° is nonempty, for then $E^\circ/2 + (-E^\circ)/2$ is a neighbourhood of the origin contained in E . Now assume $E^\circ = \emptyset$. This means that for each n , $F_n = (nE)^\circ$ is an open, dense set. Thus $\bigcap F_n$ is dense. But this is impossible, since $\bigcup nE = X$. Thus we conclude $E^\circ \neq \emptyset$. \square

For each set $E \subset X$, we let $\text{co}(E)$ denote the smallest convex subset of X containing E . This is the *convex hull* of E . For a set E , we let $\langle E \rangle$ denote the smallest subspace of X containing E , and $[E]$ denote the smallest *closed* subspace.

1.3 Bounded Linear Operators

A natural object of study in linear algebra is to understand the family of linear maps between two vector spaces X and Y . In the theory of norm spaces, we want to understand the family of *continuous* linear maps. As the next theorem shows, these linear maps are also called bounded.

Theorem 1.3. *Let X and Y be norm spaces, and $T : X \rightarrow Y$ a linear map. Then the following are equivalent:*

- T is continuous.
- T is continuous at $0 \in X$.
- T is uniformly continuous.

- There exists $M \geq 0$ such that $\|Tx\| \leq M\|x\|$ for all $x \in X$.

Proof. We leave the proof of the equivalence of (1), (2), and (3) to the reader. If T is continuous at the origin, then in particular, since $T(0) = 0$, for each $\varepsilon > 0$, there exists $\delta > 0$ such that if $\|x\| \leq \delta$, $\|Tx\| \leq \varepsilon$. But this means that $\|Tx\| \leq (\varepsilon/\delta)\|x\|$ if $\|x\| \leq \delta$, and by scaling an arbitrary x , one sees that this inequality actually holds for all $x \in X$. Conversely, if M exists such that $\|Tx\| \leq M\|x\|$, and $\varepsilon > 0$, and we set $\delta = \varepsilon/M$, then if $\|x\| \leq \delta$, $\|Tx\| \leq \varepsilon$, so T is continuous at the origin. \square

We let $B(X, Y)$ denote the class of bounded linear operators from X to Y . For each such operator, we define the *operator norm*

$$\|T\| = \sup_{x \neq 0} \frac{\|Tx\|}{\|x\|}.$$

This quantity is finite if and only if the operator is bounded by the equivalence above. One can verify that this definition gives $B(X, Y)$ the structure of a norm space, and moreover, if Y is a Banach space, then $B(X, Y)$ is a Banach space. Moreover, given $T \in B(X, Y)$ and $S \in B(Y, Z)$, we can form the composition operator $S \circ T \in B(X, Z)$, and one verifies that

$$\|S \circ T\| \leq \|S\|\|T\|.$$

Thus spaces of operators have even more algebraic and analytic structure than norm spaces. This structure is heavily exploited in the theory of *operator algebras*, covered later in these notes.

Example. If Y is not a Banach space, then $B(X, Y)$ can fail to be a Banach space. For instance, if Y is c_{00} equipped with the same norm as l_1 , then $B(l_1, Y)$ is not a Banach space. For instance, if

$$T_n(a) = \sum_{k=1}^n \frac{1}{k^2} \left(\sum_{i=k}^{\infty} a_i \right) e_k,$$

then $\|T_n - T_{n+m}\| \leq \sum_{k=n}^{\infty} (1/k)^2$, so $\{T_n\}$ is Cauchy. However, the limit of these sequence of operators is the map

$$T(a) = \sum_{k=1}^{\infty} \frac{1}{k^2} \left(\sum_{i=k}^{\infty} a_i \right) e_k$$

which is an element of $B(l_1, l_1)$, but not an element of $B(l_1, Y)$.

If $\{T_n\}$ is a sequence of operators in $B(X, Y)$ which converges in the operator norm to an operator T , then it is obvious that $Tx = \lim_{n \rightarrow \infty} T_n x$ for each $x \in X$. Pointwise convergence does not imply convergence in the operator norm, as the next example shows.

Example. Recall the Banach space c_0 , and define a sequence of operators $T_n : c_0 \rightarrow \mathbf{K}$ by letting $T_n(a) = a_n$. Then each T_n is a bounded operator, with $\|T_n\| = 1$. Thus the sequence $\{T_n\}$ does not converge to zero in the operator norm. However, for each $a \in c_0$, $\lim_{n \rightarrow \infty} T_n(a) = 0$, so the sequence $\{T_n\}$ does converge pointwise to zero.

The theory of bounded operators for finite dimensional norm spaces is effectively trivial.

Theorem 1.4. Every linear map from a finite dimensional norm space X to a norm space Y is bounded.

Proof. Let $\{e_1, \dots, e_d\}$ be a basis for X , for which, without loss of generality, we may assume $\|e_i\| = 1$ for each i . Define, for each $x \in X$ with $x = a_1 e_1 + \dots + a_d e_d$,

$$|x| = (|a_1|^2 + \dots + |a_d|^2)^{1/2}.$$

Notice that the triangle inequality implies that

$$\|x\| \leq |a_1| + \dots + |a_d| \leq \sqrt{d} (|a_1|^2 + \dots + |a_d|^2)^{1/2}.$$

Thus $\|\cdot\|$ is continuous with respect, to the topology induced by $|\cdot|$. If S is the unit sphere with respect to the norm $|\cdot|$, then the Heine-Borel theorem implies that S is compact, and in particular, $\|\cdot\|$ attains a maximum and a minimum on S . In particular, since $\|x\| \neq 0$ for each $x \in S$, this implies that there is $\varepsilon > 0$ such that $\|x\| \geq \varepsilon$ for each $x \in S$. And rescaling thus shows that $\|x\| \geq \varepsilon |x|$ for each $x \in X$. In particular, this means that $|\cdot|$ and $\|\cdot\|$ are comparable. In particular, T is continuous with respect to $\|\cdot\|$ if and only if it is continuous with respect to $|\cdot|$. And we then find that for any $x = a_1 e_1 + \dots + a_d e_d$, the triangle inequality implies

$$\|Tx\| \leq |a_1| \|Te_1\| + \dots + |a_d| \|Te_d\| \leq (|a_1| + \dots + |a_d|) \max \|Te_i\| \leq \sqrt{d} |x| \max \|Te_i\|.$$

Thus T is a continuous linear map. □

Conversely, if X is infinite dimensional, and $Y \neq 0$, then there exist discontinuous linear maps from X to Y . For instance, if we consider a basis \mathcal{B} for X containing at least countably many elements $\{e_1, e_2, \dots\}$, such that $\|e_i\| = 1$ for each i , and we find $y \neq 0$ is in Y , then we can define a linear map such that $T(e_i) = i \cdot y$, and $T(f) = 0$ for each $f \in \mathcal{B}$ with $f \neq e_i$ for each i . Then T is unbounded.

Remark. A key component of the proof above was the fact that X possessed a norm whose respective topology satisfied the Heine-Borel property. If X is infinite dimensional, no such norm can exist. This can be shown by showing that for any norm $\|\cdot\|$ on an infinite dimensional space X , and any n , there is a countable collection $\{x_1, x_2, \dots\} \subset S_X$ such that $\|x_i - x_j\| \geq 1$ for $i \neq j$. Clearly this sequence can have no accumulation points. We can define such a sequence inductively. Given $\{x_1, \dots, x_n\}$, let $X_n = \langle x_1, \dots, x_n \rangle$. Then X_n is a finite dimensional norm space, hence a Banach space, hence closed in X . In particular, if $y \notin X_n$, then $d(X_n, y) > 0$. If $\alpha = d(X_n, y)$, then $d(X_n, y/\alpha) = 1$. But this means that $\|x_i - y/\alpha\| \geq 1$ for each $i \in \{1, \dots, n\}$, so we can set $x_{n+1} = y/\alpha$.

We say a map $T : X \rightarrow Y$ between two norm spaces is an *isomorphism* if for each $x \in X$, $\|Tx\| \sim \|x\|$. In particular, this implies T is continuous and injective, but not necessarily surjective. We say X is *embedded* in Y . If T is surjective, we can define an inverse map $T^{-1} : Y \rightarrow X$, which is continuous, and we then say X and Y are isomorphic. Note that if X is a Banach space, then $T(X)$ is also a Banach space. In most other areas of math, an isomorphism has to be surjective, but unfortunately, the definition of isomorphism above has become standardized. A stronger notion than an isomorphism is an *isometry*, which is a map $T : X \rightarrow Y$ such that $\|Tx\| = \|x\|$ for all $x \in X$. Again, T need not be surjective, and we say X has been *embedded* in Y . If T is surjective, T^{-1} is also an isometry, and we say X and Y are isometrically isomorphic. Isomorphic spaces have many of the same properties, but in order to transfer all properties about norm spaces, one really needs an isometry.

Example. Let $c \subset l_\infty$ be the space of all convergent sequences $\{x_n\}$. Then c is a closed subspace, and thus a Banach space. In fact, c is isomorphic to c_0 , since we can define an operator $T : c_0 \rightarrow c$ by setting

$$T(a) = (a_0 + a_1, a_0 + a_2, \dots).$$

We have $\|T(a)\|_\infty \leq 2\|a\|_\infty$, so T is continuous. Furthermore, T is invertible, with

$$T^{-1}(b) = \left(\lim_{n \rightarrow \infty} b_n, b_0 - \lim_{n \rightarrow \infty} b_n, b_1 - \lim_{n \rightarrow \infty} b_n, \dots \right),$$

and $\|T^{-1}(b)\|_\infty \leq 2\|b\|_\infty$. Thus we conclude

$$(1/2)\|a\|_\infty \leq \|T(a)\|_\infty \leq 2\|a\|_\infty,$$

and so c and c_0 are isomorphic. Nonetheless, these spaces are not isometric. One property that c_0 has, but c does not have, is that if $a \in S_{c_0}$, then there are $b, c \in S_{c_0}$ with $b \neq c$ such that $a = (b + c)/2$. In the language of convexity, we would say that S_{c_0} has no extreme points. To see why, we define

$$b_n = \begin{cases} a_n & : |a_n| \geq 1/2 \\ a_n/2 & : |a_n| < 1/2 \end{cases} \quad \text{and} \quad c_n = \begin{cases} a_n & : |a_n| \geq 1/2 \\ 3a_n/2 & : |a_n| < 1/2 \end{cases}.$$

Clearly $a = (b + c)/2$, and $\|b\|_\infty = \|c\|_\infty = 1$. Since $\lim_{n \rightarrow \infty} a_n = 0$, eventually $|a_n| \leq 1/2$, and so $b_n \neq c_n$. On the other hand, the sequence $a = (1, 1, \dots)$ lies in S_c , yet it $b, c \in S_c$ satisfy $a = (b + c)/2$, then $b = c = a$.

A special case of linear operators is where $Y = \mathbf{K}$. The space $B(X, \mathbf{K})$ is then known as the *dual space* of X , denoted by X^* . This is one of the simplest operator spaces to analyze, and thus is a useful tool for analyzing the space X^* .

1.4 Three Fundamental Theorems

Theorem 1.5. *No Banach space X has a countably infinite basis.*

Proof. Suppose X has a countable infinite basic $\{e_1, e_2, \dots\}$. For each n , let $X_n = \langle e_1, \dots, e_n \rangle$. Then X_n is closed and has non-empty interior (if its interior was non-empty, it would be absorbing, and thus $X_n = X$). But then the Baire category theorem says that $\bigcup X_n$ is nowhere dense, which is impossible since $X = \bigcup X_n$. \square

Theorem 1.6. *Suppose $T : L^p(X) \rightarrow L^q(X)$ is an operator such that if $\{f_n\}$ is a sequence converging pointwise almost everywhere to a function f , then Tf_n converges almost every to Tf . Then T is a bounded operator.*

Proof. We apply the closed graph theorem. Suppose $\{f_n\}$ is a sequence converging to some function f in the L^p norm, and Tf_n converges to g in the L^q norm. Then $\{f_n\}$ has a subsequence $\{f_{n_k}\}$ converging almost everywhere to f , so Tf_{n_k} converges almost everywhere to Tf . But this implies that $Tf = g$, since $\{Tf_{n_k}\}$ must have a subsequence converging almost everywhere to g . \square

Theorem 1.7. *If X is a Banach space, and $Y \subset X$ is a finite dimensional subspace, then there exists a closed subspace Z of X such that $X = Y \oplus Z$.*

Proof. Let $\{y_1, \dots, y_n\}$ be a basis for Y . Hahn-Banach enables us to find $\lambda_1, \dots, \lambda_n \in X^*$ such that $\lambda_i(y_i) = 1$. Consider the continuous linear operator $\lambda : X \rightarrow \mathbb{C}^n$ given by

$$\lambda(x) = (\lambda_1(x), \dots, \lambda_n(x)).$$

Set $Z = \text{Ker}(\lambda)$. Then $X = Y \oplus Z$. \square

1.5 Quotient Spaces

Let X be a norm space, and Y a closed subspace. We can then associate a natural norm on the quotient space X/Y . We define

$$\|x + Y\| = \inf_{y \in Y} \|x + y\|.$$

If X is a Banach space, then X/Y is a Banach space. More generally, being a Banach space is a *three space property*, in the sense that if any two of $\{X, Y, X/Y\}$ are Banach spaces, then the third space is also a Banach space. The quotient map $\pi : X \rightarrow X/Y$ is continuous and surjective, since $\|x + Y\| \leq \|x\|$ for each $x \in X$. The open mapping theorem thus implies that π is an open map.

Example. *The space c_0 is a closed subspace of l_∞ , so we can consider the quotient space l_∞/c_0 . For any sequence $a \in l_\infty$, and any integer n , we have $a + c_0 = a^n + c_0$, where $a_k^n = 0$ for $k < n$, and $a_k^n = a_k$ for $k \geq n$. Thus we have*

$$\|a + c_0\| \leq \inf \|a^n\|_\infty = \limsup_{n \rightarrow \infty} |a_n|.$$

Conversely, for any sequence a , and for any sequence $b \in c_0$, we have

$$\|a + b\| = \sup |a_n + b_n| \geq \limsup |a_n + b_n| = \limsup |a_n|.$$

Thus we have $\|a + c_0\| = \limsup |a_n|$.

Theorem 1.8. *Let $T : X \rightarrow Y$ be a (not necessarily continuous) linear map, and let Z be a closed subspace of the kernel of T . Then there is a unique map $S : X/Z \rightarrow Y$ such that $T = S \circ \pi$, T is bounded if and only if S is bounded, and $\|T\| = \|S\|$.*

Proof. The existence and uniqueness of S is provided by linear algebra. If T is bounded, then

$$\|S(x + Z)\| = \|Tx\| \lesssim \|x\|,$$

Thus

$$\|S(x + Z)\| \lesssim \inf_{z \in Z} \|x + z\| = \|x + Z\|,$$

which shows S is bounded and $\|S\| \leq \|T\|$. Conversely, if S is bounded, then $T = S \circ \pi$ is bounded, and $\|T\| \leq \|S\| \|\pi\| = \|S\|$. \square

This theorem leads to a ‘first isomorphism theorem’ for Banach spaces. If $T : X \rightarrow Y$ is a bounded linear map, and K is the kernel then the induced map $S : X/K \rightarrow Y$ is an isomorphism. We conclude with an application to the theory of finite rank operators.

Theorem 1.9. *Let $T : X \rightarrow Y$ be a finite rank linear map from X to Y . Then T is bounded if and only if its kernel is closed in X .*

Proof. Suppose the kernel K of T is closed. Then X/K is a finite dimensional norm space, and so $S : X/K \rightarrow Y$ is automatically continuous, hence T is continuous. \square

If T is a linear functional, a much stronger statement can be made.

Theorem 1.10. *Let $T : X \rightarrow \mathbf{K}$ be a non-zero linear functional. Then T is bounded if and only if its kernel K is not dense in X .*

Proof. The kernel K of T has codimension one in X . If K is not dense in X , then \overline{K} is a closed, proper subspace of X , which implies $K = \overline{K}$, and thus the last theorem implies that T is continuous. \square

Another application deals with sums of closed subspaces. If M and N are closed subspaces of a norm space X , then $M + N$ need not be closed. But in one special case, the sum is closed.

Theorem 1.11. *Let M be a finite dimensional subspace of X , and let N be a closed subspace. Then $M + N$ is closed.*

Proof. Since N is closed, we can consider the quotient space X/N . For the resultant projection map $\pi : X \rightarrow X/N$, $\pi(M + N) = \pi(M)$ is a finite dimensional subspace of X/N , hence closed. By continuity, this implies that $\pi^{-1}(\pi(M)) = M + N$ is closed. \square

Chapter 2

Hilbert Spaces

2.1 Real Inner Product Spaces

Hilbert spaces are those infinite dimensional vector spaces which behave analytically as close to \mathbf{R}^n as is possible. They are therefore the spaces where we have the most powerful. Just like in a Banach space, in a Hilbert space we can measure the distances between two points via a norm. But unlike in Banach spaces, in Hilbert spaces, we can measure the angle between two vectors with an *inner product* connected to this norm.

Recall that in \mathbf{R}^2 , the angle $\theta \in [0, \pi]$ between two vectors v and w in \mathbf{R}^n satisfies the equation

$$v_1 w_1 + \cdots + v_n w_n = |v||w| \cos \theta,$$

Where $|v|^2 = v_1^2 + \cdots + v_n^2$ and $|w|^2 = w_1^2 + \cdots + w_n^2$. The left hand quantity is often easier to work with than the angle θ itself, because it is *bilinear* in the two vectors v and w . It is called the *scalar product* of the two vectors, and we denote it by $v \cdot w$. The length of a vector v can then be described via the scalar product via the equation $|v|^2 = v \cdot v$.

Inner product spaces generalize the scalar product on \mathbf{R}^n to more general settings. An inner product on a vector space V is a map associating to any two vectors v and w a scalar value $\langle v, w \rangle$, which is *bilinear*, *symmetric*, and *positive definite*, in the sense that $\langle v, v \rangle > 0$ for $v \neq 0$. A n inner product space is then a vector space equipped with a fixed inner product.

Example. For each $n \geq 0$, the vector space \mathbf{R}^n is an inner product space equipped with the scalar product $v \cdot w$.

Example. For any measure space X , the space $L^2(X)$ of square-integrable real-valued functions is an inner product space, where the inner product of $f, g \in L^2(X)$ is

$$\langle f, g \rangle = \int_{-\infty}^{\infty} f(x)g(x) \, dx.$$

The Cauchy-Schwartz inequality guarantees that this quantity is finite for any two functions. A particular example is $L^2(\mathbf{R})$, the square-integrable functions on the real line. Often more tractable for feeling out the basic theory is the space $l^2(\mathbf{Z})$ of functions $f : \mathbf{Z} \rightarrow \mathbf{R}$, such that

$$\sum_{n \in \mathbf{Z}} f(n)^2 < \infty,$$

i.e. the space of square-integrable sequences.

Any inner product space V is naturally a norm space if we define the length of a vector $v \in V$ by the equation $\|v\|^2 = \langle v, v \rangle$. One can see immediately that $\|\lambda v\| = |\lambda| \|v\|$, and that $\|v\| = 0$ if and only if $v = 0$. The triangle inequality will follow shortly, once we prove the most fundamental inequality in analysis, the Cauchy-Schwartz inequality, in the abstract setting of inner product spaces. One way to view the Cauchy-Schwartz inequality here is in terms of a generalization of the Pythagorean theorem in this setting. We say two vectors v_1, v_2 in an inner product space V are *orthogonal*, or *perpendicular*, if $\langle v_1, v_2 \rangle = 0$. This means precisely that the two vectors form a right angle with one another.

Theorem 2.1 (Pythagorean Theorem). *If $v_1, \dots, v_n \in V$ are orthogonal,*

$$\|v_1 + \dots + v_n\|^2 = \|v_1\|^2 + \dots + \|v_n\|^2$$

Proof. We just calculate that

$$\begin{aligned} \|v_1 + \dots + v_n\|^2 &= \langle v_1 + \dots + v_n, v_1 + \dots + v_n \rangle \\ &= \sum_{i,j=1}^n \langle v_i, v_j \rangle \\ &= \sum_{i=1}^n \langle v_i, v_i \rangle \\ &= \sum_{i=1}^n \|v_i\|^2 + \dots + \|v_n\|^2. \end{aligned}$$

□

Corollary 2.2. *Let V be an inner product space. Fix $v_0 \in V$. Then for any $v \in V$, the projection*

$$P_{v_0}(v) = \frac{\langle v, v_0 \rangle}{|v_0|^2} v$$

is the closest vector on the line spanned by v_0 to the vector v .

Proof. The vector $v - P_{v_0}(v)$ is orthogonal to any scalar multiple of v_0 , and so the Pythagorean theorem guarantees that for any $\lambda \in \mathbf{R}$,

$$\begin{aligned} \|v - \lambda v_0\|^2 &= \|v - P_{v_0}(v) + P_{v_0}(v) - \lambda v_0\|^2 \\ &= \|v - P_{v_0}(v)\|^2 + \|P_{v_0}(v) - \lambda v_0\|^2 \\ &\geq \|v - P_{v_0}(v)\|^2, \end{aligned}$$

with equality if and only if $\lambda v_0 = P_{v_0}(v)$. \square

Theorem 2.3 (Cauchy Schwartz Inequality). *Let V be an inner product space. Then for any $v_1, v_2 \in V$,*

$$|\langle v_1, v_2 \rangle| \leq |v_1| |v_2|.$$

Proof. The idea of this proof, in two dimensions, is intuitively simple, a trivial consequence of the Pythagorean theorem: the hypotenuse of a triangle is its longest side. We just rephrase this argument in terms of an inner product space. The two vectors $v_2 - P_{v_1}(v_2)$ and $P_{v_1}(v_2)$ are orthogonal to one another, and so the Pythagorean theorem guarantees that

$$|v_2|^2 = \|P_{v_1}(v_2)\|^2 + \|v_2 - P_{v_1}(v_2)\|^2 \geq \|P_{v_1}(v_2)\|^2 = \frac{\langle v_1, v_2 \rangle^2}{|v_1|^2}$$

Rearranging the equation and taking square roots gives the inequality. \square

Remark. The inequality in this proof shows the Cauchy-Schwartz inequality is only tight when v_1 and v_2 are close to being scalar multiples to one another, i.e. when $\|v_2 - P_{v_1}(v_2)\|$ is small. In particular, the Cauchy-Schwartz inequality is an *equality* if and only if v_1 and v_2 are scalar multiples of one another.

This theorem justifies the fact that we can define the angle $\theta \in [0, \pi]$ between two vectors by the equation

$$\langle v, w \rangle = \|v\| \|w\| \cos \theta$$

The result also implies that the norm we have defined satisfies the triangle inequality.

Theorem 2.4 (Triangle Inequality). *Let V be an inner product space. Then for any $v_1, v_2 \in V$,*

$$\|v_1 + v_2\| \leq \|v_1\| + \|v_2\|$$

Proof. We just calculate

$$\|v_1 + v_2\|^2 = \|v_1\|^2 + 2\langle v_1, v_2 \rangle + \|v_2\|^2 \leq \|v_1\|^2 + 2\|v_1\|\|v_2\| + \|v_2\|^2 = (\|v_1\| + \|v_2\|)^2$$

and then we take square roots. \square

We can now think of an inner product space as a kind of ‘infinite dimensional’ Euclidean space, with a topology given by the length function, and induced by the triangle inequality. If the metric structure is *complete*, then we say that the inner product space is a *Hilbert space*.

2.2 Complex Inner Product Spaces

Generalizing inner product spaces to vector spaces over the complex numbers presents several difficulties, because unless $V = 0$, there is no bilinear, symmetric, positive-definite form $\langle \cdot, \cdot \rangle$ on V , since we would then have for any $v \in V$,

$$\langle iv, iv \rangle = i^2 \langle v, v \rangle = -\langle v, v \rangle.$$

To obtain a definition that works in this setting, we turn to our most basic example, the real vector space \mathbf{R}^2 , which we can identify with \mathbf{C} . Here, complex conjugation plays a role in defining the distance function, removing the rotation properties of a complex number. The length of a particular complex number z is given by the equation $|z|^2 = z\bar{z}$. Over \mathbf{C}^n , we then introduce the scalar *Hermitian product*

$$(v, w) = v_1 \overline{w_1} + \cdots + v_n \overline{w_n}$$

this equation agrees with the standard scalar product when v and w have real coordinates, and also gives a notion of length, since $(v, v) > 0$ for any $v \neq 0$. Unlike the scalar product, the hermitian product is not bilinear, but instead *sesquilinear* (‘one and a half linear’ in Greek), i.e. linear in the first argument, but *antilinear* in the second, i.e. that

$$(v, w + u) = (v, w) + (v, u) \quad \text{and} \quad (v, \lambda w) = \overline{\lambda}(v, w).$$

In general, we define a *Hermitian inner product* (\cdot, \cdot) over a complex vector space to be a sesquilinear symmetric map which is positive definite. A Hermitian inner product space is just a space with a Hermitian inner product. A complex Hilbert space is just a Hermitian inner product space whose norm gives a complete metric space structure. These spaces have much the same basic theory as their real counterpart, and we will treat them identically in the sequel until we get to more advanced contexts.

Example. For any measure space X , the space $L^2(X)$ of square integrable functions is a Hermitian inner product spaces under the inner product

$$(f, g) = \int_X f(x) \overline{g(x)} dx.$$

As a special case, so is the space $l^2(\mathbf{Z})$ of square integrable sequences, equipped with the inner product

$$(f, g) = \sum_{n \in \mathbf{Z}} f(n) \overline{g(n)}$$

The conjugation generalize these spaces to complex functions.

Example. We define the Hardy space $H^2(\mathbf{D})$ to be the space of all functions $f(z)$, defined on the interior of the unit disk, and holomorphic there, such that the quantity

$$\|f\|_{H^2(\mathbf{D})} = \sup_{0 \leq r < 1} \left(\frac{1}{2\pi} \int_{-\pi}^{\pi} |f(re^{it})|^2 dt \right)^{1/2}$$

is finite. This norm is induced by the inner product

$$(f, g) = \sup_{0 \leq r < 1} \frac{1}{2\pi} \int_{-\pi}^{\pi} f(re^{it}) \overline{g(re^{it})} dt$$

which is finite by the Cauchy Schwartz inequality on $L^2[-\pi, \pi]$. Thus $H^2(\mathbf{D})$ is a Hermitian inner product space.

We claim $H^2(\mathbf{D})$ is complete, so that it is a complex Hilbert space. To do this, we will show that, as an inner product space, it is isomorphic to $l^2(\mathbf{N})$. Let us begin by noting that any $f \in H^2(\mathbf{D})$ must have a power series expansion of the form

$$f(z) = \sum_{n=0}^{\infty} a_n z^n,$$

which converges on the interior of the unit disk. Moreover, for any $0 < r < 1$, Cauchy's integral formula shows that

$$a_n = \frac{r^{-n}}{2\pi} \int_{-\pi}^{\pi} f(re^{it}) e^{-nit} dt.$$

Thus Parseval's inequality, and the orthogonality of the exponentials on any disk centred at the origin, shows that

$$\frac{1}{2\pi} \int_{-\pi}^{\pi} |f(re^{it})|^2 = \sum |a_n|^2 r^{2n}.$$

One immediate consequence of this is that the quantities

$$\frac{1}{2\pi} \int_{-\pi}^{\pi} |f(re^{it})|^2$$

are monotonically increase in r , so that

$$\|f\|_{H^2(\mathbf{D})} = \lim_{r \rightarrow 1} \left(\frac{1}{2\pi} \int_{-\pi}^{\pi} |f(re^{it})|^2 \right)^{1/2} dt.$$

Moreover, we have $\|f\|_{H^2(\mathbf{D})} = \|a\|_{l^2(\mathbf{N})}$. Thus we have found an isometry $H^2(\mathbf{D}) \rightarrow l^2(\mathbf{N})$. For any $a \in l^2(\mathbf{N})$, $\limsup_{n \rightarrow \infty} |a_n|^{1/n} \leq 1$, and so the power series

$$f(z) = \sum_{n=0}^{\infty} a_n z^n$$

is holomorphic on the interior of the unit disk, and Cauchy Schwartz shows this is clearly the inverse map to the one above. Thus $H^2(\mathbf{D})$ is a complete Hermitian inner product space.

It is a standard result of basic Fourier analysis that if $a = \{a_n\}$ lies in $l^2(\mathbf{N})$, and we consider the square integrable function

$$u(z) = \sum_{n=0}^{\infty} a_n z^n,$$

in $L^2(\mathbf{T})$, then the Poisson kernels

$$(P_r * u)(z) = \sum a_n r^n z^n$$

converge pointwise almost everywhere to u . Thus we obtain from our discussion Fatou's theorem, that for any $f \in H^2(\mathbf{D})$, and for almost every $|z| = 1$, the radial limit

$$\lim_{r \rightarrow 1} f(rz)$$

is well defined, and on the boundary defines a square-integrable function in $L^2(\mathbf{T})$. The family of all functions obtained in this way is denoted $H^2(\mathbf{T})$. It is precisely the family of functions $f \in L^2(\mathbf{T})$ such that

$$\hat{f}(n) = 0$$

for any $n < 0$.

The function $f(z) = (1-z)^{-i}$ lies in $H^2(\mathbf{T})$, but does not have a radial limit on the arc to one. If $\delta > 0$ is suitably small, it follows that for any countable set $\{w_k\}$ on \mathbf{T}

$$g(z) = \sum_{i=1}^{\infty} \delta^i f(z \overline{w_k})$$

lies in $H^2(\mathcal{D})$, but does not have radial limits on the arc towards any of the points w_k .

For any $0 < s < 1$ and $w \in \mathbf{T}$, the set $\Gamma_s(w)$ is defined to be the smallest convex set contained w and the ball of radius s at the origin. The tangential limit of a function f defined on the open unit disk at a point w is

$$\lim_{\substack{z \rightarrow w \\ z \in \Gamma_s(w)}} f(z).$$

For any $0 < s < 1$, there exists $\varepsilon > 0$ such that for $z \in \Gamma_s(w)$,

$$|(z - w) \cdot w| \geq \varepsilon \cdot (z - w).$$

Thus Cauchy Schwartz implies that if $w = e^{it_0}$, and $z \in \Gamma_s(w)$ can be written as re^{it} , then there is $C > 0$ such that

$$|t - t_0| \leq C(1 - r).$$

Thus if $f \in L^1(\mathbf{T})$, and t_0 lies in the Lebesgue set of f , then

$$|(P_r * f)(t) - f(t_0)| \leq \int |f(s + t_0) - f(t_0)| P_r(s - (t - t_0)) ds = (P_{r, t-t_0} * f)(t_0) - f(t_0),$$

where $P_{r,s_0}(s) = P_r(s - s_0)$. Now we know $\{P_r\}$ satisfies bounds of the form

$$|P_r(s)| \leq C_1/(1-r) \quad \text{and} \quad |P_r(s)| \leq C_1(1-r)/|s|^2,$$

and these bounds are sufficient to prove that $(P_r * f)(t_0)$ converges to $f(t)$ for any t_0 in the Lebesgue set of f . Similarly, where C is as above, let $P_{r,s_0}(s) = P_r(s - s_0)$ for any $|s_0| \leq C(1-r)$. We claim that we have uniform bounds of the form

$$|P_{r,s_0}(s)| \leq C_2/(1-r) \quad \text{and} \quad |P_{r,s_0}(s)| \leq C_2(1-r)/|s|^2.$$

It would then follow that as $r \rightarrow 1$, $(P_{r,s_0} * f)(t)$ converges to $f(t_0)$ for any t_0 in the Lebesgue set of f , which would complete our argument. The first bound here follows immediately from the corresponding bound for P_r . On the other hand, if $|s| \geq 2C(1-r)$, then

$$|P_{r,s_0}(s)| = |P_r(s - s_0)| \leq C_1(1-r)/|s - s_0|^2 \leq 2C_1(1-r)/|s|^2.$$

If $|s| \leq 2C(1-r)$, then

$$|P_{r,s_0}(s)| = |P_r(s - s_0)| \leq C_1/(1-r) = C_1(1-r)/(1-r)^2 \leq 4C^2C_1(1-r)/|s|^2.$$

Thus setting $C_2 = \max(2C_1, 4C^2C_1)$ implies we have shown what was required.

If W is a closed subspace of a Hilbert space H , in the sense that it is a subspace, and it is a closed set under the topology induced by the metric structure of the hermitian product, then W is also a Hilbert space, because a closed subset of a complete metric space is also complete. Conversely, a complete subspace is also closed. In the example above, the space $H^2(\mathbf{T})$ is a closed subspace of $L^2(\mathbf{T})$, and thus a Hilbert space.

Example. If Ω is an open subset of the complex plane, the space $L_a^2(\Omega)$ of holomorphic square-integrable functions is a Hilbert space. If $B_r(x)$ is a ball of radius r centered at x whose closure is contained in Ω , and f is analytic in Ω , then the mean value property guarantees that

$$f(x) = \oint_{B_r(x)} f(y) \, dy$$

It follows from Cauchy Schwarz that if $0 < r < \text{dist}(x, \partial\Omega)$, then

$$|f(x)| \leq \frac{1}{r\sqrt{\pi}} \|f\|_2$$

In particular, this implies that if a sequence of analytic functions in $L_a^2(\Omega)$ converges in L^2 norm to some function, then they must also converge locally uniformly, so the resultant function is analytic. It follows that $L_a^2(\Omega)$ is closed in $L^2(\Omega)$, and so $L_a^2(\Omega)$ is a Hilbert space. In particular, since

$$\|f\|_{L_a^2(\mathbf{D})}^2 \lesssim \|f\|_{H^2(\mathbf{D})}^2,$$

the Hardy space $H^2(\mathbf{D})$ is a closed subspace of $L_a^2(\mathbf{D})$.

There is some way to recover the intuition of the inner product measuring angles between vectors in the Hermitian setting. As we have seen in real spaces, the normalized inner product

$$\frac{\langle v, w \rangle}{|w|^2}$$

is the scalar needed to projection v onto the line spanned by the vector w . In complex vector spaces, the space spanned by a single vector is a *complex line*, or a one dimensional plane. If we set

$$\frac{\langle v, w \rangle}{|w|^2}$$

to be the value we must scale w by to obtain the projection of v onto the plane spanned by w , then we might have to rotate w , hence $\langle v, w \rangle$ might have a rotational part. Indeed, this is exactly the correct geometrical definition of the inner product, if we already have a well defined projection operation. The Pythagorean theorem, Cauchy Schwartz inequality, and triangle inequality goes through for Hermitian product spaces unperturbed, and the proofs are essentially the same, as the reader can verify. Ultimately, however, the reason that everything for real inner product spaces ‘seems’ to work for Hermitian products, is one of the main reasons we push through the initial lack of intuition.

Another approach to understanding Hermitian inner products is to use the fact that we can construct a Hermitian inner product from a real inner product. Recall that a complex vector space V is just a real vector space with a fixed twisting linear map $J : V \rightarrow V$ with $J^2 = -1$. If we already have a *real* inner product $\langle \cdot, \cdot \rangle$ on V , then our hope is to extend the real inner product to a complex inner product (\cdot, \cdot) , where the real part of the Hermitian inner product is our original real product, so that

$$\Re((v, w)) = \langle v, w \rangle$$

By geometric intuition, it is reasonable to assume that J is a twist by a right angle, so that $\langle v, Jv \rangle = 0$ for all vectors v , and that J is an isometry with respect to the inner product. It then follows that

$$\langle v, Jw \rangle = \langle Jv, J^2 w \rangle = \langle -Jv, w \rangle$$

so we cannot hope for the extension to be complex linear, only sesquilinear. If (\cdot, \cdot) is a sesquilinear form, then we calculate

$$\operatorname{Im}(v, w) = \Re(-i(v, w)) = \Re(v, iw) = \langle v, iw \rangle$$

so the *only* way to define (\cdot, \cdot) properly is as

$$(v, w) = \langle v, w \rangle + i\langle v, iw \rangle$$

The fact that (v, w) measures the projection of v onto the complex line containing w appears again, because over the real numbers, the complex line is spanned over real scalars by w and Jw , and the value above, once normalized, gives the real projection onto this two dimensional plane.

If V is any real inner product space, then we can always complexify it by embedding V in $V \oplus V$ by the trivial map $v \mapsto (v, 0)$, and defining a real inner product $\langle (v_1, v_2), (w_1, w_2) \rangle = \langle v_1, w_1 \rangle + \langle v_2, w_2 \rangle$. We can then add an artificial twisting map $J(v, w) = (-w, v)$ gives $V \oplus V$ a complex structure, and since it satisfies the properties of the discussion above, the real inner product extends to give a complex inner product. The space with all these gadgets attached is called the *complexification* of V , denoted $V^{\mathbb{C}}$. If V is a Hilbert space, then so $V^{\mathbb{C}}$ will also be a Hilbert space. In particular, because of this fact, most results for inner product spaces can be proved first for hermitian inner products, and then reduced to inner products by this embedding.

2.3 Orthogonality

Recall the notion of orthogonal, that two vectors x, y in an inner product space are orthogonal if $(x, y) = 0$. We sometimes write this statement in shorthand as $x \perp y$. Orthogonality is a symmetric relation, since if $(v, w) = 0$,

$$(w, v) = \overline{(v, w)} = \overline{0} = 0$$

The theory of right triangles essentially manifests in the theory of equations involving two to three vectors, as we have already seen in the case of the Pythagorean identity.

Theorem 2.5 (Parallelogram Law). *Let H be a Hermitian inner product space. If $x, y \in H$, then*

$$\|x + y\|^2 + \|x - y\|^2 = 2\|x\|^2 + 2\|y\|^2$$

which is geometrically intuitive, viewing $0, x, y$, and $x + y$ as vertices of a parallelogram.

Proof. We simply calculate

$$\begin{aligned} \|x + y\|^2 + \|x - y\|^2 &= \langle x + y, x + y \rangle + \langle x - y, x - y \rangle \\ &= 2\|x\|^2 + 2\|y\|^2 + \langle x, y \rangle + \langle y, x \rangle - \langle x, y \rangle - \langle y, x \rangle \\ &= 2\|x\|^2 + 2\|y\|^2 \end{aligned}$$

and this shows the theorem is true. □

2.4 Optimizing over Convex Sets

One of the biggest advantages of Hilbert spaces is that we have powerful theorems about the structure of convex sets in the space. A *convex set* in a vector space is a set C such that the line between any two points is contained in C . That is, for any $x, y \in C$, and $\lambda \in [0, 1]$, $\lambda x + (1 - \lambda)y \in C$.

Theorem 2.6. *If C is any closed convex subset of a Hilbert space, and $x \in H$ is arbitrary, then there is a unique point $y \in C$ which minimizes $\|x - y\|$, over all choices of y in the convex set.*

Proof. Since $C - x$ is a convex set, we may assume without loss of generality that $x = 0$, and so we are trying to minimize the norm over C . If $x, y \in C$ minimize this norm, with $\|x\| = \|y\| = d$, then

$$d \leq \left\| \frac{x + y}{2} \right\| \leq \frac{1}{2}(d + d) = \frac{d}{2}$$

Hence the parallelogram law implies

$$d^2 = \left\| \frac{x + y}{2} \right\|^2 = d^2 - \left\| \frac{x - y}{2} \right\|^2$$

which implies $x = y$. We know that there must a sequence $x_i \in C$ with $\|x_i\| \rightarrow d$ monotonically, and then the parallelogram law implies that if $\|x_n\|, \|x_m\| \leq d + \varepsilon$, then since $(x_n + x_m)/2 \in C$, then

$$\left\| \frac{x_n - x_m}{2} \right\|^2 = \frac{1}{2} \|x_n\|^2 + \frac{1}{2} \|x_m\|^2 - \left\| \frac{x_n + x_m}{2} \right\|^2 \leq 2d\varepsilon + \varepsilon^2 \rightarrow 0$$

hence $\{x_n\}$ is a Cauchy sequence, and therefore converges to a point x , which satisfies $\|x\| = d$. \square

Linear subspaces of a vector spaces are trivially verified to be convex sets, and we find that

Theorem 2.7. *If V is a closed linear subspace of a Hilbert space H , and $v \in V$ is the closest point to a fixed point x , then $v - x \perp V$, in the sense that $v - x \perp w$ for any $w \in V$. Conversely, if $v - x \perp V$, then v is the closest point to x .*

Proof. If $\langle v - x, w \rangle \neq 0$, then by multiplying w by a scalar value, we may assume that $\langle v - x, w \rangle < 0$, and $\|w\| = 1$, and then for $\alpha > 0$,

$$\|\alpha w + v - x\|^2 = \|v - x\|^2 + 2\alpha \Re \langle w, v - x \rangle + |\alpha|^2 \|w\|^2$$

and so as $\alpha \rightarrow 0$, we find that eventually $\|\alpha w + v - x\|^2 < \|v - x\|^2$. Conversely, if $\langle v - x, w \rangle = 0$ for all w , then for any w ,

$$\|w - x\|^2 = \|(w - v) + (v - x)\|^2 = \|w - v\|^2 + \|v - x\|^2 \geq \|v - x\|^2$$

so v minimizes the distance to x among all vectors in V . \square

For any subset X of an inner product space V , define X^\perp to be the set of vectors v such that $\langle x, v \rangle = 0$ for all $x \in X$. It is clearly a closed linear subspace of V . If V is a closed linear subspaces of a Hilbert space H , then for any $x \in H$, there is a unique $y \in V$ with $y - x \perp V$, and we may define a map $P : H \rightarrow V$ such that $Px = y$, which we call the *orthogonal projection* of H onto V .

Theorem 2.8. *If $P : H \rightarrow V$ is defined as above, then*

1. P is a linear map.
2. $\|Px\| \leq \|x\|$.

3. $P^2 = P$.

4. $\ker P = V^\perp$, and $\operatorname{im} P = V$.

Proof. If $x_0 - y_0 \perp V$ and $x_1 - y_1 \perp V$, then $(\alpha x_0 - \alpha y_0) + (\beta x_1 - \beta y_1) \perp V$, hence $P(\alpha x_0 + \beta x_1) = \alpha y_0 + \beta y_1$, and so P is a linear map. It is obvious that P is a contraction, because $\|x\|^2 = \|Px\|^2 + \|x - Px\|^2$, and $\|x - Px\|^2 \geq 0$. For any $y \in V$, $y - y \perp V$, hence $Py = y$, and so $P^2 = P$. Finally, $Py = 0$ if and only if $y \perp V$, in which case we find $y \in V^\perp$. \square

Example. Given $f \in L^2[-\pi, \pi]$, we can consider an L^2 Fourier expansion

$$f(x) = \sum_{n=-\infty}^{\infty} \hat{f}(n) e^{nit}.$$

The partial sum operators are

$$S_N f(x) = \sum_{n=-N}^N \hat{f}(n) e^{nit}$$

and these operators are precisely the projection operators onto the finite dimensional subspace of $L^2[-\pi, \pi]$ spanned by $\{e_{-N}, \dots, e_N\}$. We also have an orthogonal projection

$$Pf(x) = \sum_{n=0}^{\infty} \hat{F}(n) e^{nit} = \frac{1}{2\pi i} \int_{\partial \mathbf{D}} \frac{f(w)}{w - z} dw.$$

whose image we can identify with $H^2(\mathbf{T})$, e.g. the space of all functions which form the boundary values of holomorphic functions in $H^2(\mathbf{D})$.

Corollary 2.9. For any closed subspace V of a Hilbert space H , $(V^\perp)^\perp = V$. If X is an arbitrary subset of a Hilbert space H , then $(X^\perp)^\perp = \overline{\operatorname{span}}(X)$ is the closure of the linear span of the elements of X .

Proof. If P is the orthogonal projection of H onto V , and Q is the orthogonal projection of H onto V^\perp , then $Px + Qx = x$ for all $x \in H$, and so the kernel of Q is exactly the set of vectors x for which $Px = x$, from which we conclude that $x \in V$. If X is arbitrary, then $\overline{\operatorname{span}}(X)$ is contained in $(X^\perp)^\perp$, because $(X^\perp)^\perp$ is a closed subspace of H containing all elements in X .

Conversely, we conclude that $X^\perp = \overline{\text{span}}(X)^\perp = X^\perp$, because if $\langle y, x \rangle = 0$ for all $x \in X$, then surely $\langle y, \sum \alpha_i x_i \rangle = \sum \alpha_i \langle y, x_i \rangle = 0$, so $y \perp \text{span}(X)$, and if $v_i \rightarrow v$, with $v_i \in \text{span}(X)$, then by continuity of the inner product we find that $0 = \langle y, v_i \rangle \rightarrow \langle y, v \rangle$, hence $\langle y, v \rangle = 0$, and so $y \perp \overline{\text{span}}(X)$. \square

Corollary 2.10. *A subspace V is dense in H if and only if $V^\perp = 0$.*

Proof. For then the closure of V is $(V^\perp)^\perp = 0^\perp = H$. \square

Given two Hilbert spaces H_0 and H_1 , we can define a canonical Hilbert space structure on $H_0 \oplus H_1$ by letting $\langle x_0 + x_1, y_0 + y_1 \rangle = \langle x_0, y_0 \rangle + \langle x_1, y_1 \rangle$. The induced topology is precisely the product topology, because $x_i + y_i \rightarrow x + y$ if and only if $x_i \rightarrow x$, and $y_i \rightarrow y$. The orthogonal projection of H onto V is incredibly important to the theory of Hilbert spaces, because it gives us a direct sum decomposition of H into $V \oplus V^\perp$. The Pythagorean theorem tells us that for any $v \in V$ and $w \in V^\perp$,

$$\|v + w\|^2 = \|v\|^2 + \|w\|^2$$

and so the map $x \mapsto Px + (1 - P)x$ is actually an *isometry* between vector spaces (it preserves the inner product). This is the basic notion of equivalence between various Hilbert spaces, and so we really can think of H as being composed of the two closed subspaces V and V^\perp .

2.5 The Riesz Representation Theorem

Many problems in functional analysis can be reduced to the analysis of certain linear functions between vector spaces. The easiest linear functions to analyze are the *linear functionals*, which are linear maps $f : V \rightarrow K$, where K is the base field of the vector space V . We will say a linear function $f : V \rightarrow K$ is *bounded* if there exists $M > 0$ such that $|f(x)| \leq M\|x\|$ for all $x \in V$. The bounded linear maps are exactly the continuous linear maps (and to verify continuity, we need only verify continuity at zero, so bounded linear maps are essentially ‘uniformly continuous’). If f and g are continuous, then $\lambda f + \gamma g$ is continuous, and so the set V^* of bounded linear functionals on V is a vector space, called the *dual space* of V . On the dual space of a norm space, we may define a dual norm by letting $\|f\|$ be the smallest number M such that $|f(x)| \leq M\|x\|$ holds for all x . Since

$|f(x) + g(x)| \leq |f(x)| + |g(x)|$, we find $\|f + g\| \leq \|f\| + \|g\|$, and it is easy to verify that $\|\lambda f\| = |\lambda| \|f\|$.

Given a Hilbert space H , we can use the inner product structure on H to transport vectors in H to functionals in H^* . That is, for a given $x \in H$, we define $x^* \in H^*$ to be the functional $x^*(y) = \langle x, y \rangle$. The Cauchy-Schwarz inequality tells us that $\|x^*\| = \|x\|$, so the dual map is an isometry. The Riesz representation theorem tells us that this isometry is actually invertible.

Theorem 2.11 (Riesz). *For any bounded linear function $f \in H^*$, there exists a unique vector $x \in H$ such that $f(y) = \langle x, y \rangle$ for all $y \in H$.*

Proof. Assume without loss of generality that $f \neq 0$. Let V be the kernel of f . Then V is a closed linear subspace of H , and V^\perp is one dimensional by the first isomorphism theorem. Consider the orthogonal projection P onto V^\perp . If x is a nonzero vector in V^\perp , with $\|x\| = 1$ and $f(x) = \gamma$, then for each y there is $\lambda(y)$ for which $f(y) = f(Py) = \lambda(y)$, and what's more, $\langle y, x \rangle = \langle Py, x \rangle = \lambda(y)$, and so we have found a unique y as specified, and y is unique because the dual map is an isometry. \square

Example. *For any measure space X , $L^2(X)$ is a Hilbert space, and so for any bounded linear functional $\alpha : L^2(X)^*$, there is a unique $g \in X$ such*

$$\alpha(f) = \int_X f(x) \overline{g(x)} \, dx.$$

Thus studying linear functionals on $L^2(X)$ reduces to integration theory.

2.6 Constructing an Orthonormal Basis

In the theory of general infinite dimension vector spaces, the standard definition of a basis (which in the sequel we will call a *Hamel basis*, expanding a vector in a finite linear combination of basis elements, does not necessarily lead to a useful theory because of analytical problems. In the setting of Banach spaces, we can get around this by instead considering infinite linear combinations of vectors. A set $\{e_n\}$ is a *Schauder basis* for a Banach space X such that for any $x \in X$, there is a unique family of scalars λ_n such that $x = \sum_n \lambda_n e_n$, where $\sum_n \lambda_n e_n$ is defined as the limit of the net of finite

sums corresponding to the sequence. In the sequel, by a basis, we will mean a Schauder basis.

If $\{e_1, \dots, e_n\}$ are an *orthonormal* family of vectors in an inner product space V , in the sense that they are pairwise orthogonal, and $\|e_n\| = 1$ for all i , then the Pythagorean identity implies

$$\left\| \sum_n a_n e_n \right\| = \left(\sum |a_n|^2 \right)^{1/2}.$$

Using slightly more analysis, and using the fact that we are working in a Hilbert space H , this result remains true for an infinite family of orthonormal vectors $\{e_n\}$, provided that $\sum |a_n|^2 < \infty$, because then the partial sums of $\sum a_n e_n$ are Cauchy, and because a Hilbert space is *complete*, the partial sums must converge to an infinite sum. For any vector v , the sum $\sum_n (v, e_n) e_n$ also converges to some vector in H , and by continuity of the inner product, the vector $v - \sum (v, e_n) e_n$ is orthogonal to all of the vectors $\{e_n\}$. Thus we find that

$$\|v\|^2 = \|v - \sum (v, e_n) e_n\|^2 + \sum_n |(v, e_n)|^2 \geq \sum_n |(v, e_n)|^2.$$

This inequality is *Bessel's inequality*. If $v = \sum (v, e_n) e_n$, then we find that

$$\|v\|^2 = \sum_{i=1}^{\infty} |(v, e_n)|^2.$$

This equation is known as *Parseval's identity*. We now give a series of conditions on $\{e_n\}$ guaranteeing that for any $v \in H$, $v = \sum_n (v, e_n) e_n$.

Theorem 2.12. *The following properties of an orthonormal family $\{e_n\}$ are equivalent:*

- (a) *For any vector $v \in H$, $v = \sum_n (v, e_n) e_n$.*
- (b) *Finite linear combinations of the e_i are dense in H .*
- (c) *If $(v, e_n) = 0$ for all n , then $v = 0$.*
- (d) *Parseval's identity holds.*

We then say the sequence is an orthonormal basis for the space.

Proof. (a) clearly implies (b). Assuming (b), note that if $(v, e_i) = 0$ for all i , then in particular, $(v, w) = 0$ for all w which can be expressed as finite linear combinations of the e_i . But we may choose w_n with $w_n \rightarrow v$ in L^2 , and then since $|(v, w_n - v)| \leq \|v\| \|w_n - v\| \rightarrow 0$, we conclude that $0 = (v, w_n) = (v, v) + (v, w_n - v)$ converges to (v, v) , so $(v, v) = 0$, and hence $v = 0$. If v is given, then Bessel's inequality ensures that

$$\sum_{i=1}^{\infty} |(v, e_i)|^2 \leq \|v\|^2 < \infty$$

and so the partial sums of the sequence $\sum_{i=1}^{\infty} |(v, e_i)|^2$ are Cauchy, and because we are working in a Hilbert space, converge to some vector w . Now $v - w$ is orthogonal to e_i for each i , since

$$(w, e_i) = \lim_{n \rightarrow \infty} \sum_{j=1}^n (v, e_j)(e_j, e_i) = (v, e_i)$$

and therefore $v = w$, proving (a). (a) implies (d) is trivial, and if (d) holds, then Pythagoras' identity implies that

$$\|v - \sum_{i=1}^n (v, e_i) e_i\|^2 = \sum_{i=n+1}^{\infty} |(v, e_i)|^2 \rightarrow 0$$

so (a) holds. □

A topological space is *separable* if it has a countable dense subset. If H is a separable Hilbert space, then by performing the Gram Schmidt process to a countable dense subset, one finds an orthogonal sequence, which is easily verified to be an orthonormal basis for the Hilbert space. Thus *every separable Hilbert space has a orthonormal basis*. With some more work on the geometry of Hilbert spaces, and the axiom of choice, we can show that an arbitrary Hilbert space has an orthonormal basis, if we allow bases ranging over an uncountable index set. Before we do this, however, we consider some applications.

Example. *Fourier analysis was the first area of mathematics where the ideas of orthogonality were employed in infinite dimensions. Since $[-\pi, \pi]$ is a finite*

measure space, the space $L^2[-\pi, \pi]$ is contained within $L^1[-\pi, \pi]$, and so if $f \in L^2[-\pi, \pi]$, we can consider its Fourier series

$$\hat{f}(n) = \frac{1}{2\pi} \int_{-\pi}^{\pi} f(x) e^{-nix} dx.$$

If we define $e_n(x) = e^{nix}$, then $\hat{f}(n) = (f, (2\pi)^{-1}e_n)$, and

$$(e_n, e_m) = \int_{-\pi}^{\pi} e^{(n-m)ix} = \begin{cases} 2\pi & n = m \\ 0 & n \neq m \end{cases}$$

It follows that $\{(2\pi)^{-1}e_n\}$ is an orthonormal sequence in $L^2[-\pi, \pi]$. Using the basic fact from Fourier analysis that for any $f \in C[-\pi, \pi]$ with $f(-\pi) = f(\pi)$, the partial sums of

$$\sum \hat{f}(n) e^{nix}$$

converge uniformly to f , and the fact that continuous functions are dense in $L^2[-\pi, \pi]$, we conclude that the exponentials $\{(2\pi)^{-1}e_n\}$ are an orthonormal basis for $L^2[-\pi, \pi]$. As a result, we obtain the classic Parseval's inequality

$$\int_{-\pi}^{\pi} |f(x)|^2 dx = \sum_{n=-\infty}^{\infty} |\hat{f}(n)|^2,$$

and that for any $f \in L^2[-\pi, \pi]$, the partial sums of

$$\sum \hat{f}(n) e^{nix}$$

converge in $L^2[-\pi, \pi]$ to f .

Example. Viewing \mathbf{T} as the boundary of the unit disk in \mathbf{C} , equipped with the normalized Lebesgue measure. Then for appropriately smooth $f, g \in L^2(\mathbf{T})$, we can interpret (f, g) as a contour integral, i.e.

$$(f, g) = \int_{\mathbf{T}} \frac{f(z) \overline{g(z)}}{2\pi i z} dz,$$

where the torus is given the standard orientation. Now consider the Möbius transformation

$$z = \frac{i - x}{i + x},$$

which is an oriented mapping from \mathbf{R} to $\mathbf{T} - \{-1\}$. Now

$$dz = -2i \frac{1}{(x+i)^2}.$$

Thus

$$(f, g) = \frac{1}{\pi} \int_{-\infty}^{\infty} f\left(\frac{i-x}{i+x}\right) \overline{g\left(\frac{i-x}{i+x}\right)} \frac{1}{1+x^2} dx.$$

Thus if we consider the map $T : L^2(\mathbf{T}) \rightarrow L^2(\mathbf{R})$ given by

$$Tf(x) = f\left(\frac{i-x}{i+x}\right) \frac{\pi^{-1/2}}{x+i},$$

then T is an isometry, i.e. $(Tf, Tg) = (f, g)$ for each $f, g \in L^2(\mathbf{T})$. In particular, since $L^2(\mathbf{T})$ has an orthonormal basis of the form $\{z^n : n \in \mathbf{Z}\}$, we conclude that $L^2(\mathbf{R})$ has an orthonormal basis of the form

$$\pi^{-1/2} \left(\frac{i-x}{i+x}\right)^n \frac{1}{i+x}.$$

We can generalize to non separable spaces by an application of Zorn's lemma. If $\{S_\alpha\}$ is a chain of orthonormal sets, then $\bigcup S_\alpha$ is an orthonormal set, so we may apply Zorn's lemma to conclude that a maximal orthonormal set S exists. We claim this is an orthonormal basis, which is equivalent to showing that $S^\perp = \{0\}$. And indeed, if S^\perp was nontrivial, we could find $x \in S^\perp$ with $\|x\| = 1$, and then $S \cup \{x\}$ would be an orthonormal set, contradicting maximality. Thus orthonormal bases exist on any Hilbert space.

Theorem 2.13. *Any two orthonormal bases of a Hilbert space H have equal cardinality.*

Proof. If S is an orthonormal basis for a Hilbert space H , then one of two possibilities hold:

- Some orthonormal basis of H is finite: Then H is finite dimensional, and the result follows by standard linear algebra.
- There does not exist a finite orthonormal basis for H : Let T be another orthonormal basis. For each $x \in H$, the set $T_x = \{e \in T : (e, x) \neq 0\}$ is at most countable, since the sum $\sum |(e, x)|^2$ converges. But $T = \bigcup_{x \in S} T_x$ since S is an orthonormal basis, which means that

$$\#(T) \leq \#(S) \cdot \#(\mathbf{N}) = \#(S).$$

□

We may therefore define the *dimension* of a Hilbert space to be the cardinality of any orthonormal basis. This actually uniquely defines the Hilbert space, for if H_0 and H_1 are Hilbert spaces with orthonormal bases $\{e_i\}$ and $\{f_i\}$, then we have an isometry $T : H_0 \rightarrow H_1$ given by

$$T(\sum a_i e_i) = \sum a_i f_i.$$

This is well defined, for the sum $\sum a_x x$ converges if and only if $\sum |a_x|^2 < \infty$, and the same criterion shows $\sum a_x f(x)$ converges.

Example. For any set X equipped with the counting measure, we have an orthonormal basis on $l^2(X)$, given by the family of functions $\{e_x : x \in X\}$, where $e_x(y) = \mathbf{I}(x = y)$. The family is certainly orthonormal. If $f \in l^2(X)$ satisfies $f(x) = \langle f, e_x \rangle = 0$ for all $x \in X$, then $f = 0$, so the set is a basis. Thus the dimension of the Hilbert space $l^2(X)$ is $\#(X)$. If X is not countable, then $l^2(X)$ is not separable.

Example. Define an inner product on the space $AP_0(\mathbf{R})$ of all finite linear combinations of exponentials on \mathbf{R} , i.e. an inner product on the span of $\{e^{2\pi i \lambda x} : \lambda \in \mathbf{R}\}$, by declaring the set to be orthonormal. For any two $f, g \in AP_0(\mathbf{R})$, one can verify that

$$\langle f, g \rangle = \lim_{T \rightarrow \infty} \frac{1}{2T} \int_{-T}^T f(t) \overline{g(t)} dt,$$

i.e. this is an inner product space in which each of the exponentials in $\{e^{2\pi i \lambda x}\}$ are separable. Write the resulting norm as $\|f\|_{AP(\mathbf{R})}$. The completion of $AP_0(\mathbf{R})$ is denoted $AP(\mathbf{R})$. Note that since $AP(\mathbf{R})$ has an uncountable basis, it is not separable.

A continuous function $f \in C_b(\mathbf{R})$ is called uniformly almost periodic if it is the uniform limit of a family of functions in $AP_0(\mathbf{R})$. By virtue of the fact that $\|f\|_{AP(\mathbf{R})} \leq \|f\|_{L^\infty(\mathbf{R})}$, almost periodic functions can be identified with a certain subfamily of elements of $AP(\mathbf{R})$. If f_0 is uniformly almost periodic, and $0 < \varepsilon \leq 1$, then we can find $f \in AP_0(\mathbf{R})$ such that $\|f - f_0\|_{L^\infty(\mathbf{R})} \leq \varepsilon$. The function f is the linear combination of finitely many exponentials oscillating at a certain family of frequencies $\{\lambda_1, \dots, \lambda_n\}$, i.e. we can write $f = \sum a_i e_{\lambda_i}$. If

$|\lambda_1|, \dots, |\lambda_n| \leq \alpha$, then

$$\begin{aligned}
|f(t+s) - f(t)| &\leq \sum a_i |e_{\lambda_i}(t+s) - e_{\lambda_i}(t)| \\
&\leq \left(\sum |a_i|^2 \right)^{1/2} \left(\sum |\lambda_i s + \mathbf{Z}|^2 \right)^{1/2} \\
&= \|f\|_{AP(\mathbf{R})} \left(\sum |\lambda_i s + \mathbf{Z}|^2 \right)^{1/2} \\
&\leq 2\|f_0\|_{L^\infty(\mathbf{R})} \left(\sum |\lambda_i s + \mathbf{Z}|^2 \right)^{1/2}.
\end{aligned}$$

Consider the sequence $X(n)_i = \lambda_i n + \mathbf{Z}$ in \mathbf{T}^n . Then

$$|X(n)| = \left(\sum |\lambda_i n + \mathbf{Z}|^2 \right)^{1/2}.$$

By Weyl equidistribution, this sequence is either periodic, or equidistributed on some finite union of cosets of a subtorus of \mathbf{T}^n . Fix some $R > 0$ large enough such that $X(1), \dots, X(R)$ is a ε cover of this finite union of cosets. Then, given any interval $I \subset \mathbf{R}$ of length $R + 1$, we can find an integer $n_0 \in I$ such that $[n_0 - R, n_0] \subset I$. Find $n \in [1, R]$ such that $|X(n) - X(n_0)| \leq \varepsilon$, then $n - n_0 \in I$, and $|X(n - n_0)| = |X(n) - X(n_0)| \leq \varepsilon$. Thus

$$|f(t+n) - f(t)| \leq |f(t+n) - f(t)| + \varepsilon \leq (1 + 2\|f_0\|_{L^\infty(\mathbf{R})}) \cdot \varepsilon.$$

Adjusting parameters, we conclude that for any almost periodic function f_0 , and any $\varepsilon > 0$, we can find $L > 0$ such that any interval I of length L contains an ‘almost period’ x_0 , such that $|f_0(x+x_0) - f_0(x)| \leq \varepsilon$ for all $x \in \mathbf{R}$. Conversely, suppose that for each $\varepsilon > 0$, we can find $L > 0$ such that any interval of length L contains a ε almost period for f . Consider a trigonometric polynomial g , periodic on $[-10L, 10L]$, such that $\|f - g\|_{L^\infty[-10L, 10L]} \leq \varepsilon$. Then for any $x_1 \in \mathbf{R}$, we can find $x_0 \in \mathbf{R}$ with $|x_1 - x_0| \leq L$ such that $|f(x+x_0) - f(x)| \leq \varepsilon$ for all $x \in \mathbf{R}$. Then

$$f(x_1) = f(x_1 - x_0) + O(\varepsilon) = g(x_1 - x_0) + O(\varepsilon) = g(x_1 - x_0 + 10L)$$

$$|f(x+x_0) - f(x)| \leq \varepsilon.$$

We can then find a trigonometric polynomial $g = \sum a_n e_{n/x_0}$ such that $\|f - g\|_{L^\infty[0, x_0]} \leq \varepsilon$. **TODO FINISH THIS PROOF PLUS LOOK UP BESICOVITCH, ALMOST PERIODIC FUNCTIONS.**

Chapter 3

Basics of Operator Theory

3.1 Operators on Hilbert Spaces

If H is a Hilbert space, then the characterization of H^* is very easy - it is naturally isomorphic to H itself. And Hilbert spaces form very neat isomorphism classes, so, roughly speaking, the theory of problems about the elements of Hilbert spaces is roughly speaking, a solved problem. Modern research on Hilbert spaces is really about the bounded operators on a Hilbert space, which form a rich and nontrivial class for which many open problems remain unsolved. We begin the study of such operators in this chapter. Given two Hilbert spaces H_1 and H_2 , we let $B(H_1, H_2)$ denote the bounded linear maps from H_1 to H_2 .

The first powerful tool in the study of operators on Hilbert spaces is the existence of *adjoints*. Given a linear map $T : X \rightarrow Y$ between Banach spaces, we obtain a natural adjoint map $T^* : Y^* \rightarrow X^*$. If X and Y are Hilbert spaces, then X is naturally isomorphic to X^* , and Y is naturally isomorphic to Y^* . Thus we can actually define the adjoint as a map $T^* : Y \rightarrow X$. Let us introduce the adjoint more precisely.

Lemma 3.1. *If (\cdot, \cdot) is a continuous sesquilinear form on a Hilbert space H , then there is a unique $T \in B(H)$ such that $(x, y) = \langle x, Ty \rangle$ for all $x, y \in H$.*

Proof. Since (\cdot, \cdot) is continuous, the uniform boundedness principle implies that there exists a quantity $L > 0$ such that for any $x, y \in H$,

$$|(x, y)| \leq L \|x\| \|y\|.$$

For each x , (\cdot, x) is a bounded linear functional, so there is a unique $Tx \in H$ such that $(y, x) = \langle y, Tx \rangle$. Now

$$\begin{aligned}\langle y, T(ax_1 + bx_2) \rangle &= (y, ax_1 + bx_2) \\ &= \bar{a}(y, x_1) + \bar{b}(y, x_2) \\ &= \bar{a}\langle y, Tx_1 \rangle + \bar{b}\langle y, Tx_2 \rangle \\ &= \langle y, aTx_1 + bTx_2 \rangle,\end{aligned}$$

it follows that T is a linear map, i.e. $T(ax_1 + bx_2) = aTx_1 + bTx_2$. Moreover, since

$$|\langle x, Ty \rangle| = |(x, y)| \leq L\|x\|\|y\|,$$

it follows that $\|Ty\| \leq L\|y\|$.

$$\langle x, My \rangle = (x, y) \leq \|(\cdot, \cdot)\|\|x\|\|y\|$$

So $\|M\| \leq \|(\cdot, \cdot)\|$, but also

$$(x, y) = \langle x, My \rangle \leq \|M\|\|x\|\|y\|$$

so $\|(\cdot, \cdot)\| \leq \|M\|$, and we have equality of norms. \square

Theorem 3.2. *Let H_1 and H_2 be Hilbert spaces. If $T : H_1 \rightarrow H_2$ is a bounded linear map, then there is a unique bounded linear map $T^* : H_2 \rightarrow H_1$ such that*

$$\langle Tx, y \rangle = \langle x, T^*y \rangle$$

The map T^ is the adjoint of T . We have $T^{**} = T$, and $\|T^*\| = \|T\|$.*

Proof. The map $(x, y) = \langle Tx, y \rangle$ is a sesquilinear form, so there is a unique T^* such that $(x, y) = \langle x, T^*y \rangle$. By taking suprema over $x \in H$ with $\|x\| = 1$ on the right hand side of the equation defining T^* , we have

$$\|T^*y\| = \sup_{\|x\|=1} \langle Tx, y \rangle \leq \|T\|\|y\|.$$

Thus $\|T^*\| \leq \|T\|$. It is simple to see that $T^{**} = T$, because for any $x \in H_1$ and $y \in H_2$,

$$\langle T^{**}x, y \rangle = \langle x, T^*y \rangle = \langle Tx, y \rangle.$$

But this means by symmetry that $\|T\| = \|T^{**}\| \leq \|T^*\|$, which implies $\|T\| = \|T^*\|$. \square

Example. If we are working over the field \mathbf{R} , then every operator $B(\mathbf{R}^n, \mathbf{R}^m)$ can be identified with a unique matrix in $M_{n,m}(\mathbf{R})$, and the adjoint corresponds to taking the transpose of a matrix. Working over the field \mathbf{C} , every operator $B(\mathbf{C}^n, \mathbf{C}^m)$ can be identified with a unique matrix in $M_{n,m}(\mathbf{C})$, and the adjoint here is the conjugate transpose.

Example. Given $\phi \in L^\infty(X)$, where X is a σ -finite measure space, we have an operator $M_\phi \in B(L^2(X))$ defined by $M_\phi(f) = \phi f$. Then

$$\langle M_\phi f, g \rangle = \int \phi f \bar{g} = \int f \overline{g\bar{\phi}} = \langle f, \bar{\phi}g \rangle$$

Thus $M_\phi^* = M_{\bar{\phi}}$.

The adjoint operation $*$ is an antilinear operator from $B(H_1, H_2)$ to $B(H_2, H_1)$, such that $(ST)^* = T^*S^*$. This is verified by calculation, since

$$\langle STx, y \rangle = \langle Tx, S^*y \rangle = \langle x, T^*S^*y \rangle$$

We also have $\|T^*T\| = \|T\|^2$, since

$$\langle T^*Tx, x \rangle = \langle Tx, Tx \rangle = \|T\|^2 \|x\|^2$$

so $\|T^*T\| \leq \|T\|^2$, and this inequality is attained in the suprema, since if $x_i \in B_H$ is such that $\|Tx_i\| \rightarrow \|T\|$, then

$$\langle T^*Tx_i, x_i \rangle = \|Tx_i\|^2 \rightarrow \|T\|^2$$

Thus $*$ is what is known as an *involution*, and makes $B(H)$ into a C^* algebra.

Next, we prove a result which generalizes the rank-nullity theory theorem from finite dimensional linear algebra.

Theorem 3.3. Fix a linear operator $T : H \rightarrow H$ on a Hilbert space H . Then

$$H = \text{Im}(T) \oplus \text{Ker}(T^*) = \text{Ker}(T) \oplus \text{Im}(T^*).$$

Proof. By symmetry, it suffices to prove the first equality. The kernel of T^* is closed, and so it suffices to show that $\text{Ker}(T^*) = \text{Im}(T)^\perp$. If $x \in \text{Im}(T)$ and $y \in \text{Ker}(T^*)$, then we can write $x = Tx'$, and then

$$\langle x, y \rangle = \langle Tx', y \rangle = \langle x', T^*y \rangle = \langle x', 0 \rangle = 0.$$

Thus $\text{Ker}(T^*) \subset \text{Im}(T)^\perp$. Conversely, if $y \in \text{Im}(T)^*$, then for any $x \in H$,

$$\langle T^*y, x \rangle = \langle y, Tx \rangle = 0.$$

This means $T^*y = 0$. □

In the special case where T is *self adjoint*, i.e. $T^* = T$, we conclude that

$$H = \text{Ker}(T) \oplus \text{Im}(T).$$

We will later see this remains true for the class of normal operators, and moreover, we can refine this orthogonal sum when T is compact to give a complete eigenfunction expansion.

Any operator $M : H \rightarrow H$ can be expressed as $M = T + iS$, where T and S are self adjoint operators, by defining

$$T = \frac{1}{2}(M + M^*) \quad S = \frac{1}{2i}(M - M^*)$$

That T and S are unique follows because if T_1 and T_2 are self-adjoint, and $T_1 = iT_2$, then we also have

$$T_1 = T_1^* = (iT_2)^* = -iT_2^* = -iT_2$$

Thus $iT_2 = -iT_2$, which means $T_1 = T_2 = 0$. We denote T by $\text{Re}(M)$ and $S = \text{Im}(M)$. An operator $M : H \rightarrow H$ is *normal* if $M^*M = MM^*$, or equivalently, if $\text{Re}(M)\text{Im}(M) = \text{Im}(M)\text{Re}(M)$.

Theorem 3.4. *If T is a self-adjoint operator, then*

$$\|T\| = \sup\{\langle Tx, x \rangle : \|x\| \leq 1\}$$

Proof. Let M be the supremum. Clearly $M \leq \|T\|$. If $\|x\|, \|y\| \leq 1$,

$$\begin{aligned} \langle T(x \pm y), x \pm y \rangle &= \langle Tx, x \rangle \pm \langle Ty, x \rangle \pm \langle Tx, y \rangle + \langle Ty, y \rangle \\ &= \langle Tx, x \rangle \pm \overline{\langle Tx, y \rangle} \pm \langle Tx, y \rangle + \langle Ty, y \rangle \\ &= \langle Tx, x \rangle \pm 2\Re\langle Tx, y \rangle + \langle Ty, y \rangle \end{aligned}$$

Then, subtracting, we find

$$4\Re\langle Tx, y \rangle = \langle T(x + y), x + y \rangle - \langle T(x - y), x - y \rangle$$

Thus

$$4\Re\langle Tx, y \rangle \leq M(\|x + y\|^2 + \|x - y\|^2) = 2M(\|x\|^2 + \|y\|^2) \leq 4M$$

Let λ be such that $\langle Tx, y \rangle = \lambda|\langle Tx, y \rangle|$. If we replace x by λx , we find

$$|\langle Tx, y \rangle| = \bar{\lambda}\langle Tx, y \rangle = \langle T(\bar{\lambda}x), y \rangle \leq M$$

Hence

$$\|Tx\| = \sup\{|\langle Tx, y \rangle| : y \in B_H\} \leq M$$

and therefore $\|T\| \leq M$. □

Remark. We will later see that this result remains true for all *normal operators* T , using the spectral theory of such operators. The main idea is that the spectral theory implies the existence of an *approximate eigenvalue* $\lambda \in \mathbf{C}$ for T with $|\lambda| = \|T\|$, i.e. a scalar such that for any $\varepsilon > 0$, there is $x \in H$ with $\|x\| = 1$ such that $\|Tx - \lambda x\| \leq \varepsilon$. This implies that if M is as in the proof above, then

$$M \geq |\langle Tx, x \rangle| \geq |\lambda| - \varepsilon = \|T\| - \varepsilon,$$

and we can take $\varepsilon \rightarrow 0$.

Corollary 3.5. *Let $T : H \rightarrow H$ be self adjoint. Then*

$$\langle Tx, x \rangle = 0$$

for all $x \in H$ if and only if $T = 0$.

Corollary 3.6. *If $\langle Tx, x \rangle = \langle Sx, x \rangle$ for each $x \in H$, then $T = S$.*

Corollary 3.7. *An operator $T : H \rightarrow H$ is normal if and only if for any $x \in H$,*

$$\|Tx\| = \|T^*x\|.$$

Proof. If T is normal, then for any $x \in H$, we calculate that

$$\|Tx\|^2 = \langle Tx, Tx \rangle = \langle T^*Tx, x \rangle = \langle TT^*x, x \rangle = \langle T^*x, T^*x \rangle = \|T^*x\|^2.$$

Conversely if $\|Tx\| = \|T^*x\|$ for all $x \in H$, then we find for any $x \in H$,

$$\langle T^*Tx, x \rangle = \langle TT^*x, x \rangle.$$

Since T^*T and TT^* are both self-adjoint, this means that $T^*T = TT^*$, i.e. T is normal. \square

The next few lemmas show that these are the unique self-adjoint operators with this property. This is essentially like beginning with \mathbf{C} , and identifying \mathbf{R} as those elements invariant under complex conjugation.

Theorem 3.8. *T is normal if and only if $\|Tx\| = \|T^*x\|$ for each x .*

Proof. For then, for each x ,

$$\|Tx\|^2 = \langle Tx, Tx \rangle = \langle T^*Tx, x \rangle$$

$$\|T^*x\|^2 = \langle T^*x, T^*x \rangle = \langle TT^*x, x \rangle$$

which implies $T^*T = TT^*$. The converse repeats backwards through the proof, since the norm of $\|T\|$ is specified by its action through the inner product. \square

Theorem 3.9. *Any normal operator T has the following properties.*

- (a) $\ker(T) = \ker(T^*)$.
- (b) $T(H)$ is dense in H if and only if T is injective.
- (c) T is a surjective isomorphism if and only if $\|Tx\| \geq C\|x\|$ for some C .
- (d) If $Tx = \lambda x$, then $T^*x = \bar{\lambda}x$.
- (e) If λ and γ are distinct eigenvalues, then the eigenspaces are orthogonal to one another.

Proof. To prove (a) notice that $\|Tx\| = 0$ if and only if $\|T^*x\| = 0$. (b) follows since

$$T(H)^\perp = \ker(T^*) = \ker(T)$$

If $\|Tx\| \geq C\|x\|$, then $T(H)$ is closed, and dense, therefore $T(H) = H$. The converse follows by definition. If T is normal, then $\lambda - T$ is also normal, since $(\lambda - T)^* = \bar{\lambda} - T^*$, so $\ker(\lambda - T) = \ker(\lambda - \bar{\lambda})$. Finally, if $Tx = \lambda x$, and $Ty = \gamma y$, then

$$\lambda \langle x, y \rangle = \langle Tx, y \rangle = \langle x, T^*y \rangle = \langle x, \bar{\gamma}y \rangle = \gamma \langle x, y \rangle$$

so if $\lambda \neq \gamma$, then $\langle x, y \rangle = 0$. \square

Theorem 3.10. *If $U \in B(H)$, then the following are equivalent.*

- (a) U is unitary.
- (b) $GL(H) = H$ and $\langle Ux, Uy \rangle = \langle x, y \rangle$.
- (c) $GL(H) = H$ and $\|Ux\| = \|x\|$, so U is an isometry from H to itself.

If any of these results are true, we say U is a unitary operator.

Proof. If U is unitary, then U is invertible, so $GL(H) = H$, and

$$\langle Ux, Uy \rangle = \langle U^* Ux, y \rangle = \langle x, y \rangle$$

If (b) holds, then

$$\|Ux\|^2 = \langle Ux, Ux \rangle = \langle x, x \rangle = \|x\|^2$$

If (c) holds, then

$$\langle U^* Ux, x \rangle = \langle Ux, Ux \rangle = \|Ux\|^2 = \|x\|^2 = \langle x, x \rangle$$

Thus $U^* U$ is the identity. Since U is injective and surjective, the open mapping theorem tells us U is invertible, so $U^* = U^{-1}$. \square

This shows that unitary operators are the natural automorphisms in the category of Hilbert spaces.

Lemma 3.11. *An operator $T \in B(H)$ in a complex Hilbert space is self adjoint if and only if $\langle Tx, x \rangle \in \mathbf{R}$ for each $x \in H$.*

Proof. If T is self adjoint, then

$$\langle Tx, x \rangle = \langle x, Tx \rangle = \overline{\langle Tx, x \rangle}$$

Conversely, if $\langle Tx, x \rangle \in \mathbf{R}$ for each $x \in H$, then

$$\langle Tx, x \rangle = \overline{\langle x, Tx \rangle} = \langle x, Tx \rangle = \langle T^* x, x \rangle$$

for each x , from which we conclude $T = T^*$. \square

Theorem 3.12. *If P is a projection with image W , then the following are equivalent:*

- (a) P is self-adjoint.
- (b) P is normal.
- (c) $W = \ker(P)^\perp$.
- (d) $\langle Px, x \rangle = \|Px\|^2$.

In any of the three cases, we say P is the orthogonal projection onto W , and W uniquely defines P .

Proof. If P is self-adjoint, then P is trivially normal. If P is normal, then

$$P(H) = \ker(P^*)^\perp = \ker(P)^\perp$$

If (c) holds, then every $x = y + Px$, where $Py = 0$. Then

$$\langle P(y + Px), y + Px \rangle = \langle Px, y + Px \rangle = \langle Px, Px \rangle = \|Px\|^2$$

If (d) holds, then $\langle Px, x \rangle \in \mathbf{R}$ for each x , so P is self-adjoint. If (c) holds, then $\langle Px, x \rangle$ □

If H is a Hilbert space, then two orthogonal projections P_1 and P_2 might not necessarily commute. As a basic example, let us consider only projections onto a line. Fix $\|v_1\| = \|v_2\| = 1$, and consider $P_i x = \langle x, v_i \rangle v_i$. Then

$$P_1 P_2 x = \langle x, v_2 \rangle \langle v_2, v_1 \rangle v_1$$

and

$$P_2 P_1 x = \langle x, v_1 \rangle \langle v_1, v_2 \rangle v_2.$$

If the lines are orthogonal, i.e. $\langle v_1, v_2 \rangle = 0$, then $P_1 P_2 = P_2 P_1 = 0$ trivially commute. On the other hand, if the lines are not orthogonal, then we find that v_1 and v_2 are scalar multiples of one another, in which case $P_1 = P_2$. Thus one dimensional projections only commute if they are equal to one another, or projections onto orthogonal subspaces.

Theorem 3.13. *Let H be a Hilbert space, and let P_1 and P_2 be orthogonal projections onto subspaces W_1 and W_2 of H . Let $U_1 = (W_1 \cap W_2)^\perp \cap W_1$ and $U_2 = (W_1 \cap W_2)^\perp \cap W_2$. Then the following are equivalent:*

- P_1 and P_2 commute.
- $P_1 P_2$ is an orthogonal projection.
- U_1 is perpendicular to U_2 .

Proof. Suppose P_1 and P_2 commute, and let $P = P_1 P_2 = P_2 P_1$. Then

$$P^2 = P_1 P_2 P_1 P_2 = P_1^2 P_2^2 = P_1 P_2 = P.$$

Thus P is a projection. Moreover,

$$P^* = (P_1 P_2)^* = P_2 P_1 = P_1 P_2 = P.$$

Thus P is an orthogonal projection. Since $P = P_1 P_2$, the image of P is contained in W_1 . Since $P = P_2 P_1$, the image of P is contained in W_2 . Thus the image of P is contained in $W_1 \cap W_2$. Conversely, if $x \in W_1 \cap W_2$, then $Px = x$. Thus P is projection onto $W_1 \cap W_2$. On the other hand, suppose $P_1 P_2$ is an orthogonal projection. Then $P_1 P_2$ is self adjoint, so $P_1 P_2 = (P_1 P_2)^* = P_2 P_1$, so P_1 and P_2 commute.

Now suppose P_1 and P_2 commute. Write $P_i = P_0 + P_{U_i}$, where P_0 is orthogonal projection onto $W_1 \cap W_2$, and \tilde{U}_i is orthogonal projection onto U_i . Then

$$P_1 P_2 = (P_0 + P_{U_1})(P_0 + P_{U_2}) = P_0 + P_{U_1} P_{U_2}$$

and

$$P_2 P_1 = (P_0 + P_{U_2})(P_0 + P_{U_1}) = P_0 + P_{U_2} P_{U_1}.$$

If P_1 and P_2 commute, $P_1 P_2 = P_2 P_1 = P_0$, and so $P_{U_1} P_{U_2} = P_{U_2} P_{U_1} = 0$. Thus U_1 is orthogonal to U_2 . Conversely, if U_1 is orthogonal to U_2 , then $P_{U_1} P_{U_2} = P_{U_2} P_{U_1} = 0$, and from the equations above we see P_1 and P_2 commute. \square

3.2 Compact Operators

We already know operators on an infinite dimensional Banach space are more tricky than their finite dimensional counterparts. Thus we restrict our attention to operators which do not ‘spread themselves out’ too much. An operator is *compact* if the image of any bounded set is precompact. It suffices to verify that the image of the unit ball is precompact. Given two Banach spaces X and Y , we shall let $K(X, Y)$ denote the space of all compact operators from X to Y .

Theorem 3.14. $K(X, Y)$ is a closed subspace of $B(X, Y)$, such that for any $M \in K(X, Y)$, $N \in B(Y, Z)$, and $L \in B(Z', X)$, $NML \in K(Z', Z)$.

Proof. Suppose M and N are compact. Let U be a bounded subset of X . Then

$$(\lambda M + \gamma N)(U) = \lambda M(U) + \gamma N(U)$$

And the sum of two precompact sets is precompact. Thus $\lambda M + \gamma N$ is precompact. If U is a bounded subset of Z' , then $L(U)$ is a bounded subset of X , so $ML(U)$ is a precompact subset of Y , implying $NML(U)$ is a precompact subset of Z , for if every sequence y_i has a convergent subsequence, so must $N(y_i)$.

To prove $K(X, Y)$ is closed we rely on the fact that a subset of a complete metric space is precompact if and only if, for any $\varepsilon > 0$, the subset can be covered by finitely many balls of radius ε . Now suppose we have $M_1, M_2, \dots \in K(X, Y)$, and $M_i \rightarrow N$. We claim N is compact. Let U be a bounded subset of X , such that $\|x\| < K$ for each $x \in U$. Then for each $\varepsilon > 0$, there are finitely many balls of radius ε which cover $M_i(U)$, for each i . If $\|M_i - N\| < \varepsilon/K$, then $\|(M_i - N)(x)\| < \varepsilon$ for each $x \in U$, so that if $M_i x$ is contained in some ball with center y of radius ε , then Nx is contained in the ball with center y and radius 2ε . Thus $N(U)$ can be covered with finitely many balls of radius 2ε , and is therefore precompact. \square

Corollary 3.15. *$K(X)$ is a closed, two-sided ideal of $B(X)$.*

Remark. It is a theorem of Calkin that $K(H)$ is the only closed, two-sided ideal of $B(H)$ for any Hilbert space H . Thus we are lead to consider the *Calkin Algebra* $B(H)/K(H)$, which is simple, and quite useful in K -theory.

The most basic examples of compact operators are bounded *finite rank* operators, i.e. operators whose range is finite dimensional. Indeed, bounded operators map bounded sets to bounded sets, and a closed bounded set in a finite dimensional space is compact. Because $K(X, Y)$ is closed, any limit of a family of finite rank operators is compact. Conversely, the family of all finite-rank operators is dense in $K(H)$ for any Hilbert space H (we say H possesses the *approximation property*).

Theorem 3.16. *Let H be a Hilbert space. Then $T : H \rightarrow H$ is compact if and only if T is the limit in the operator norm of a sequence of finite rank operators $\{T_n\}$.*

Proof. Suppose T is compact. Then the image H_0 of T must be separable. Let $\{e_n\}$ be a basis for H_0 , and let $T_{\leq N} = P_{\leq N}T$, where P_N is the projection onto the first N elements of the basis for H_0 . If it was not true that $T_{\leq N}$ converged to T in the operator norm, we could find $\varepsilon > 0$, and two sequences $\{N_i\}$ and $\{x_i\}$ such that $\|x_i\| = 1$ for all i , and $\|P_{\leq N_i}Tx_i - Tx_i\| \geq \varepsilon$.

Since T is compact, by thinning the subsequence we may assume that $Tx_i \rightarrow y$. Then

$$\|P_{\leq N_i}y - y\| \geq \|P_{\leq N_i}Tx_i - Tx_i\| - \|P_{\leq N_i}(Tx_i - y)\| - \|Tx_i - y\| \geq \varepsilon - 2\|Tx_i - y\|.$$

Now $P_{\leq N_i}y \rightarrow y$, so

$$0 = \limsup_{i \rightarrow \infty} \|P_{\leq N_i}y - y\| \geq \limsup_{i \rightarrow \infty} \varepsilon - 2\|Tx_i - y\| = \varepsilon,$$

which gives a contradiction. Thus $T_{\leq N}$ converges to T . \square

Remark. Finding a Banach space *without* the approximation property (i.e. finding a compact operator which cannot be approximated by finite rank operators) proved to be an incredibly difficult problem in 20th century research on functional analysis. Such an operator was only found in 1972.

Corollary 3.17. *Let H be a Hilbert space. If T has finite rank, then T^* has finite rank. A bounded operator T on H is compact if and only if T^* is compact.*

Proof. Suppose T has finite rank. Then we can find an orthogonal basis $\{e_n\}$ for H such that

$$Tx = \sum_{i,j=1}^M a_{ij} \langle x, e_i \rangle e_j = \sum_{i,j=1}^M a_{ij} S_{ij}x.$$

But

$$\langle S_{ij}x, y \rangle = \langle x, e_i \rangle \langle e_j, y \rangle = \langle x, \langle y, e_j \rangle e_i \rangle = \langle x, S_{ji}y \rangle.$$

Thus $S_{ij}^* = S_{ji}$, and so by linearity,

$$T^*y = \sum_{i,j=1}^M \overline{a_{ij}} S_{ji}y.$$

Thus T^* has finite rank.

Now if T is a compact operator on H , then $T = \lim T_n$, where T_n is a finite rank operator. Then $T^* = \lim T_n^*$ is the limit of finite rank operators, so T^* is also compact. The reverse follows by symmetry. \square

Example. The simplest operators on a Hilbert space are the diagonal operators, i.e. those operators $T : H \rightarrow H$ given in terms of an orthogonal basis $\{e_n\}$ and $\{\lambda_n\}$ such that $Te_n = \lambda_n e_n$. Then

$$\|T\| = \sup_n |\lambda_n|,$$

If $S(\varepsilon) = \{n : |\lambda_n| \geq \varepsilon\}$ is infinite for any ε , then $\{e_n : n \in S\}$ is a bounded subset of H , but $Te_n = \lambda_n e_n$ does not have a convergent subsequence, and so T is not a compact operator. It follows that if T is compact, $\lambda_n \rightarrow 0$ as $n \rightarrow \infty$. Conversely, if $\lambda_n \rightarrow 0$ as $n \rightarrow \infty$, then the finite rank operators

$$T_{\leq N}x = \sum_{n=1}^N \lambda_n \langle x, e_n \rangle e_n$$

converge in operator norm to T as $N \rightarrow \infty$, and so T is a compact operator. Thus a bounded diagonalizable operator on a separable Hilbert space is compact precisely when its eigenvalues converge to zero. More generally, for a diagonal operator on a non separable Hilbert space, every subsequence of eigenvalues must converge to zero, which implies all but countably many eigenvalues are zero.

A useful way to verify compactness in spaces of continuous functions is to use the Arzela-Ascoli criterion. A family of functions $\{g_\alpha : X \rightarrow \mathbf{C}\}$ on a topological space X is *equicontinuous* at $x \in X$ if, for any $\varepsilon > 0$, there is an open set U containing x such that $\|g_\alpha(a) - g_\alpha(b)\| < \varepsilon$ for any $a, b \in U$ and any index α .

Theorem 3.18 (Arzela-Ascoli). *If X is a compact metric space, the precompact subsets of $C(X)$ are precisely those subsets which are bounded and equicontinuous.*

Proof. TODO □

Thus if X is a Banach space, and Y is a compact metric space, then a bounded operator $T : X \rightarrow C(Y)$ is compact precisely when T maps bounded subsets of $C(X)$ to equicontinuous subsets of $C(Y)$.

Example. The best examples of compact operators are given by integral operators with an appropriately regular kernel. Given two compact Hausdorff

spaces X and Y equipped with finite Borel measures μ and ν , and a continuous function $K : X \times Y \rightarrow \mathbf{C}$, we can define an operator by the integral formula

$$Tf(y) = \int_X K(x, y) f(x) dx.$$

Then for any $f \in L^1(X)$, $Tf \in C(Y)$, which indicates a gain of regularity and a possible compactness of the operator T . We claim T is compact from $L^p(X)$ to $L^q(Y)$ for any $1 \leq p, q \leq \infty$. To see this, we note that any integral operator with kernel

$$\tilde{K}(x, y) = \sum a_{ij} f_i \otimes g_j$$

for some continuous functions $f_i \in C(X)$ and $g_j \in C(Y)$, has finite rank and is therefore compact. The Stone-Weirstrass theorem implies such functions are dense in $C(X \times Y)$. Thus we can approximate any $K \in C(X \times Y)$ by a kernel \tilde{K} of the form above. But for any $1 \leq p, q \leq \infty$ we have

$$\|T\|_{L^p \rightarrow L^q} \lesssim \|K\|_{L^\infty(X \times Y)},$$

so it follows that T is a limit in the operator norm of finite rank operators, and thus compact.

Example. A Hilbert-Schmidt operator between two Hilbert spaces H_1 and H_2 is a bounded operator $T : H_1 \rightarrow H_2$ such that in two orthogonal bases $\{e_i\}$ and $\{f_j\}$,

$$\|T\|_{HS} = \left(\sum_{i,j} |\langle Te_i, f_j \rangle|^2 \right)^{1/2} = \left(\sum_i \|Te_i\|_{H_2}^2 \right)^{1/2} < \infty.$$

This norm is independant of the choice of bases, as can be verified by a simple expansion. If T is Hilbert Schmidt, so is T^* , and $\|T^*\|_{HS} = \|T\|_{HS}$. Moreover, we have a bound $\|T\| \leq \|T\|_{HS}$. Any Hilbert-Schmidt operator is compact, because if $P_{\leq N}$ is projection onto the first N vectors in a basis $\{e_k\}$, then $P_{\leq N}T$ has finite rank, and $P_{\leq N}T \rightarrow T$ as $N \rightarrow \infty$ in the operator norm, because it converges in the Hilbert-Schmidt norm. A basic example of a Hilbert-Schmidt operator is an integral operator T between two measure spaces X and Y with kernel $K \in L^2(X \times Y)$, in which case $\|T\|_{HS} = \|K\|_{L^2(X \times Y)}$. To see this, we let $\{e_i\}$ and $\{f_j\}$ be orthogonal basis on $L^2(X)$ and $L^2(Y)$. Since $\bar{e}_i \otimes f_j$ is then an orthogonal basis for $L^2(X \times Y)$, we calculate that

$$\|T\|_{HS}^2 = \sum |\langle Te_i, f_j \rangle|^2 = \sum |\langle K, \bar{e}_i \otimes f_j \rangle|^2 = \|K\|_{L^2(X \times Y)}^2.$$

The space of all Hilbert-Schmidt operators between two Hilbert spaces H_1 and H_2 forms a Hilbert space itself, $B_{HS}(H_1, H_2)$, where the inner product is

$$\langle T_1, T_2 \rangle_{HS} = \sum \langle T_1 e_i, f_j \rangle \overline{\langle T_2 e_i, f_j \rangle}.$$

One can also consider slightly different variants of Hilbert-Schmidt operators. For instance, the family of Hilbert-Schmidt operators between two measure spaces X and Y are the family of all operators $T : L^2(X) \rightarrow L^2(Y)$ induced by $K \in L^2(X \times Y)$ by setting

$$Tf(x) = \int K(x, y)f(y) dy.$$

If $\{e_n\}$ is an orthogonal basis for $L^2(X)$, and $\{e_m\}$ is an orthogonal basis for $L^2(Y)$, then $e_n \otimes e_m$ is a basis for $L^2(X \times Y)$, and so

$$\|K\|_{L^2(X \times Y)}^2 = \sum |\langle K, e_n \otimes e_m \rangle|^2 = \sum_{n,m} |\langle Te_n, e_m \rangle|^2.$$

Thus a Hilbert-Schmidt integral operator is a Hilbert Schmidt operator of the kind above. Conversely, if T is a normal, compact operator from $L^2(X)$ to itself, then there exists an orthonormal basis $\{e_n\}$ diagonalizing T with eigenvalues $\{\lambda_n\}$. If $K(x, y) = \sum_n \lambda_n e_n(x) e_n(y)$, then K is a Hilbert-Schmidt kernel which gives the operator T .

Example. Let B be the unit ball in \mathbf{R}^d , and consider a kernel $K : B \times B \rightarrow \mathbf{C}$ such that $|K(x, y)| \lesssim |x - y|^{-(d-s)}$, for some $s > 0$. Then

$$\|K\|_{L_y^\infty(B)L_x^1(B)} \quad \text{and} \quad \|K\|_{L_x^\infty(B)L_y^1(B)} < \infty,$$

so Schur's Lemma implies that if T is the integral operator corresponding to K , then $\|Tf\|_{L^p(\mathbf{R}^d)} \lesssim \|f\|_{L^p(\mathbf{R}^d)}$ for $1 \leq p \leq \infty$. For $s > d/2$, this operator is a Hilbert-Schmidt operator, but not for $s \leq d/2$. Nonetheless, for any $s > 0$, T is a compact operator. To see this, for each n , define T_n to be the integral operator corresponding to the kernel $K_n(x, y) = \mathbf{I}(|x - y| \geq 1/n)K(x, y)$. Then K_n is a Hilbert-Schmidt kernel, and so every operator in the family $\{T_n\}$ is compact. Moreover, we have

$$\lim_{n \rightarrow \infty} \|K - K_n\|_{L_y^\infty L_x^1} = \lim_{n \rightarrow \infty} \|K - K_n\|_{L_x^\infty L_y^1} = 0,$$

and so $T_n \rightarrow T$ in the operator norm. Thus T is compact.

TODO: Trace class operators.

TODO: If T is compact, T^* is compact.

3.3 Spectral Theory for Compact Operators

The spectral theory of operators on finite dimensional vector spaces shows that any self-adjoint matrix can be diagonalized. Here we try and generalize this result to infinite dimensional spaces. In general, this is not so simple to generalize. The most basic reformulation (any self-adjoint operator on a Hilbert space has a basis of eigenvectors) is *not correct*, because of the existence of *almost eigenvectors*.

Example. Let $T : L^2[0,1] \rightarrow L^2[0,1]$ be the multiplication operator given by $Tf(x) = x \cdot f(x)$. Then T has no eigenvectors, i.e. $Tf = \lambda f$ precisely when $f = 0$. On the other hand, for any $\lambda \in [0,1]$, we can find a sequence of ‘almost eigenvectors for λ ’, i.e. a sequence of functions $f_n \in L^2[0,1]$ such that $\|f_n\|_{L^2[0,1]} = 1$, and

$$\lim_{n \rightarrow \infty} \|Tf_n - \lambda f_n\|_{L^2[0,1]} = 0.$$

There are several alternate generalizations which *do* work for an arbitrary self-adjoint operator on a Hilbert space (in fact, for any normal operator), which are best treated in the more advanced theory of operators. For now, we limit ourselves to studying operators on a Hilbert space which are ‘almost finite dimensional’, i.e. they are compact. It is simple to see why compactness is useful in this regime: if a compact operator $T : H \rightarrow H$ has an sequence $\{x_n\}$ such that $\|x_n\| = 1$ and $\lim_{n \rightarrow \infty} \|Tx_n - \lambda x_n\| = 0$ for some $\lambda \neq 0$, then the compactness of T shows that Tx_n/λ must have a convergent subsequence, converging to some $y \in H$ with $\|y\| = 1$, and then continuity shows that

$$\|Ty - \lambda y\| = \lim_{n \rightarrow \infty} \|T^2x_n/\lambda - Tx_n\| \leq \lambda^{-1} \|T\| \lim_{n \rightarrow \infty} \|Tx_n - \lambda x_n\| = 0.$$

Thus $Ty = \lambda y$, a sequence of almost eigenvectors implies the existence of an *exact eigenvector*. For a normal operator T , we let $\sigma(T)$ be the set of *almost eigenvalues* for T , and $\sigma_p(T)$ the set of *exact eigenvalues*. We have shown that if T is compact and normal, then $\sigma(T) = \sigma_p(T)$.

To begin our study, let us introduce the basic notation. If X is a Banach space, and $T : X \rightarrow X$ is bounded, for each λ , we let $X_\lambda = \{x : Tx = \lambda x\}$. Then X_λ is a closed subspace of X . The scalar λ is called an *eigenvalue* if $X_\lambda \neq \{0\}$, and any nonzero $x \in X_\lambda$ is an eigenvector with eigenvalue λ . It turns out that in infinite dimensions, one not only has the phenomena of

eigenvalues, but also *almost* eigenvalues, which makes the family of eigenvalues not so powerful a description of operators. Thus we replace this description with the *spectrum* of a bounded operator $T : X \rightarrow X$, which is the set $\sigma(T)$ of all scalars λ such that $T - \lambda$ is *not* an isomorphism between X and itself. If T is normal, then $T - \lambda$ is an isomorphism if and only if an inequality of the form

$$\|(T - \lambda)x\| \gtrsim \|x\|$$

holds for all $x \in H$. In particular, in this case $\sigma(T)$ consists of all λ such that for any $\varepsilon > 0$, we can find $x \in H$ such that $\|Tx - \lambda x\| \leq \varepsilon \|x\|$, i.e. such that we can find x such that $Tx \approx \lambda x$, up to arbitrary precision. Thus $\sigma(T)$ detects not only eigenvalues, but also *almost eigenvalues*.

Example. A self-adjoint operator need not have any eigenvalues. For instance, if $T : L^2[0,1] \rightarrow L^2[0,1]$ is given by $Tf(x) = xf(x)$, then T is self-adjoint, but has no eigenvalues, for if $Tf = \lambda f$, for any λ , then $xf(x) = \lambda f(x)$, which implies $f = 0$. On the other hand, T any $\lambda \in [0,1]$ is almost an eigenvalue, in the following sense. If $f(x) = \mathbf{I}(x \in I)$, where I is an interval of length δ containing λ , then $|Tf(x) - \lambda f(x)| \leq \delta \mathbf{I}(x \in I)$, and so

$$\|Tf - \lambda f\|_{L^2[0,1]} \leq \delta |I|^{1/2} \leq \delta \|f\|_{L^2[0,1]}.$$

Thus if δ is small, then we have $Tf \approx \lambda f$, and we can take δ as small as we like, for any $\lambda \in [0,1]$. Thus this calculation shows $T - \lambda$ is not invertible for any $\lambda \in [0,1]$. Since $T - \lambda$ does have an inverse for $\lambda \notin [0,1]$, we have actually calculated that $\sigma(T) = [0,1]$. Thus the spectrum detects not only eigenvalues, but also almost eigenvalues.

Example. As another example of this phenomenon, let $\{\lambda_n\}$ be an enumeration of all rational numbers in $[0,1]$. Thus the bounded operator $T : l^2(\mathbf{N}) \rightarrow l^2(\mathbf{N})$ given by

$$(Ta)(n) = \lambda_n a(n)$$

Then $Te_n = \lambda_n e_n$, so T has eigenvectors corresponding to eigenvalues for all rational numbers in $[0,1]$. On the other hand, $T - \lambda$ is not invertible for any $\lambda \in [0,1]$, so $\sigma(T - \lambda) = [0,1]$.

The main result of this section is that a compact operator on a Hilbert space H can always be diagonalized, with a sequence of eigenvalues λ_i converging to zero, and such that H can be decomposed into the orthogonal sum of the eigenspaces H_{λ_i} , and for $i \neq 0$, H_{λ_i} is finite dimensional.

Theorem 3.19. *If X is a Banach space, Y is a closed subspace, T is a compact operator on X , and $\lambda \neq 0$ is a scalar such that*

$$\inf_{\substack{y \in Y \\ \|y\|=1}} (T - \lambda)y = 0,$$

then $Y \cap X_\lambda \neq 0$.

Proof. Let y_i be a sequence in Y with $\|y_i\| = 1$ such that $(T - \lambda)y_i \rightarrow 0$. Since T is a compact operator, we may assume without loss of generality that Ty_i converges to some z . But then, by continuity,

$$\lim_i \lambda y_i = \lim_i \lambda y_i - Ty_i + Ty_i = \lim_i y_i - Ty_i + \lim_i Ty_i = 0 + z$$

Since $\lambda \neq 0$, this means $y_i \rightarrow \lambda^{-1}z$, which means $\|z\| = \lambda$, and in particular, $z \neq 0$. But then $z \in Y$, and $z = Mz$, which completes the proof since $z \in Y \cap X_\lambda$. \square

TODO: IF λ is an eigenvalue for T , it is an eigenvalue for T^* .

Corollary 3.20. *Let $T : X \rightarrow X$ be a compact operator. If $\lambda \in \sigma(T)$ is non-zero, then $\dim(X_\lambda) < \infty$, $(T - \lambda)X$ is closed and $\text{codim}(T - \lambda) < \infty$, and there is n such that*

$$\ker[(T - \lambda)^n] = \ker[(T - \lambda)^{n+1}]$$

and

$$(T - \lambda)^n X = (T - \lambda)^{n+1} X$$

Proof. First note that $T = \lambda$ restricted to X_λ . Since X_λ is closed, this implies the identity operator is compact on X_λ , hence X_λ is finite dimensional. Since X_λ is finite dimensional, we can find a subspace Y of X such that $X = X_\lambda \oplus Y$. $T - \lambda$ restricted to Y is bounded and injective, and the last theorem implies, since $Y \cap X_\lambda = (0)$, that there is $C > 0$ such that

$$\|(T - \lambda)y\| \geq C\|y\|$$

for $y \in Y$. Thus $T - \lambda$ is an isomorphism from Y to a subspace Z of X , which is therefore complete and thus closed. Note that

$$Z^\perp = [(T - \lambda)X]^\perp = \{\phi \in X^* : \phi \circ (T - \lambda) = 0\} = \ker(T^* - \lambda).$$

Since T is compact, so is T^* . Because λ must also be an eigenvalue of T^* , the previous part of this argument therefore implies that

$$\text{codim}(T - \lambda) = \dim[Z^\perp] < \infty.$$

TODO □

Now we restrict to the study of compact self-adjoint operators on a Hilbert space. It follows from the following theorem that if T is a compact self-adjoint operator on a separable Hilbert space, then $\sigma(T)$ is precisely the set of eigenvalues for T .

Lemma 3.21. *If H is a Hilbert space, $T : H \rightarrow H$ is bounded and self-adjoint, and $\lambda_1 \neq \lambda_2$, then H_{λ_1} and H_{λ_2} are orthogonal.*

Proof. We begin by assuming without loss of generality that λ_1 and λ_2 are real, since a self-adjoint operator can only have real eigenvalues. If $x_1 \in H_{\lambda_1}$ and $x_2 \in H_{\lambda_2}$, then

$$\lambda_1 \langle x_1, x_2 \rangle = \langle Tx_1, x_2 \rangle = \langle x_1, Tx_2 \rangle = \lambda_2 \langle x_1, x_2 \rangle,$$

so since $\lambda_1 \neq \lambda_2$, $\langle x_1, x_2 \rangle = 0$. □

Theorem 3.22. *If $T : H \rightarrow H$ is a compact self-adjoint operator on a separable Hilbert space, then H has an orthonormal basis $\{e_i\}$ such that $Te_i = \lambda_i e_i$, where $\{\lambda_i\}$ have no limit point but zero.*

Proof. Assume without loss of generality that $T \neq 0$. Since

$$\|T\| = \sup \langle Tx, x \rangle$$

we can find a sequence x_i with $\|x_i\| = 1$, such that

$$\langle Tx_i, x_i \rangle \rightarrow \lambda,$$

where $\lambda = \pm \|T\|$. Since T is compact, we may assume $Tx_i \rightarrow y$. Then

$$\langle y, x_i \rangle \rightarrow \|T\|,$$

so $y \neq 0$. But

$$\|Tx_i - \lambda x_i\|^2 = \|Tx_i\|^2 + \lambda^2 - 2\lambda \text{Re}(\langle Tx_i, x_i \rangle) \rightarrow 0.$$

Thus $\lambda x_i \rightarrow y$, which means $Ty = \lim T(\lambda x_i) = \lambda y$, so an eigenvector exists. A Zorn's lemma argument then allows us to find a orthogonal basis of eigenvectors for T . □

Remark. Let $\{T_\alpha\}$ be a family of *commuting* self-adjoint compact operators on a Hilbert space H . Then we claim that H has an orthogonal basis simultaneously diagonalizing the operators T_α . We can do this by an application of the well-ordering principle. Suppose we have a basis $\{e_i\}$ which simultaneously diagonalizes $\{T_\alpha\}$ for $\alpha < \alpha_0$. Then the eigenspaces $\bigcap_{\alpha < \alpha_0} H_{\alpha, \lambda_\alpha}$ are T_{α_0} -invariant for all choices of scalars λ_α , since if $x \in \bigcap_{\alpha < \alpha_0} H_{\alpha, \lambda_\alpha}$,

$$T_\alpha(T_{\alpha_0}x) = T_{\alpha_0}(T_\alpha x) = \lambda_\alpha T_{\alpha_0}x.$$

Because each of these spaces is T_{α_0} -invariant, we can diagonalize each of these spaces, finding an orthonormal basis on each one, and then put them back together, and we will have found a basis which simultaneously diagonalizes $\{T_\alpha : \alpha \leq \alpha_0\}$. This means there is no smallest index α_0 such that we cannot simultaneously diagonalize $\{T_\alpha : \alpha \leq \alpha_0\}$, which means we can simultaneously diagonalize all of the operators.

3.4 Strong and Weak Convergence

It is often a useful tool to weaken the notion of convergence of operators. We say a sequence of operators $\{T_n : X \rightarrow Y\}$ *converges in norm* to an operator $T : X \rightarrow Y$ if $\|T - T_n\| \rightarrow 0$. We say $\{T_n\}$ *strongly converges* to T if $T_n x$ converges in norm of Tx for every $x \in X$ (note that this is, somewhat confusingly, *weaker* than norm convergence). Finally, $\{T_n\}$ *weakly converges* to T if $T_n x$ converges to Tx in the weak topology for every $x \in X$.

Example. *The only operators which are the norm limit of a family of finite rank operators are compact operators. On the other hand, every operator on a Hilbert space is a strong limit of a sequence of finite rank operators. TODO WHAT ABOUT BANACH SPACE?*

Chapter 4

Vector Valued Integration

Ubiquitous to analysis is the integral, defined to be $\int f(x) dx$, defined to be an average of the infinitesimal values of the function f across a certain domain. In this context, f is almost always real or complex valued. But there are many spaces in which an average can be taken. Indeed, in order to define the average of a set of values a_1, \dots, a_N , given by

$$\frac{a_1 + \dots + a_N}{N}$$

we need only that we can add the a_n , and multiply them by a scalar. Thus the natural place to generalize the integration of functions is to those which take their values in a topological vector space. In this chapter we analyze such a theory.

The weakest definition of the integral exploits the fact that the integral is linear. Because integration is the limit of finite sums, and linear functions are interchangeable with summation, for any continuous linear operator $\Lambda : X \rightarrow Y$, and function $f : \Omega \rightarrow X$, we ought to have

$$\Lambda \left(\int_{\Omega} f \right) = \int_{\Omega} (\Lambda \circ f)$$

In particular, if Λ is a linear functional, the right hand side is just a scalar-valued function, and we know how to integrate those. In fact, we can *define* the *weak integral* $\int_{\Omega} f$ such that this equation holds for any continuous linear functional Λ . In particular, if X^* separates X , then such an integral is unique, and the only question of applicability is when such an

integral exists. The conditions of the next theorem, which show that such an integral exists, are automatically satisfied if X is a Fréchet space.

Theorem 4.1. *If X is a topological vector space such that X^* separates points, and μ is a Borel probability measure on a compact Hausdorff space Ω , and $f : \Omega \rightarrow X$ is a continuous function such that the closed convex hull of $f(\Omega)$ is compact, then the weak integral exists.*

Proof. Let H be the closed convex hull of $f(\Omega)$. Given a finite family of linear functionals $L = \{\Lambda_1, \dots, \Lambda_N\}$, we consider the family E_L of all points $x \in H$ such that

$$\Lambda_n(x) = \int_{\Omega} (\Lambda_n \circ f)(x) d\mu$$

If we can prove that each E_L is non-empty, then by the compactness of H , the intersection of all such E_L will be nonempty, and this will be the required weak integral. If we define a vector y by

$$y_n = \int_{\Omega} (\Lambda_n \circ f) d\mu$$

then it suffices to show that the image of H under the map $L(x) = (\Lambda_1 x, \dots, \Lambda_N x)$ contains y . Because L is a linear map into \mathbf{C}^n , $L(H)$ is compact and convex. Thus if y was not an element of this set, we can separate y from H by a linear functional. In other words, there are constants c_1, \dots, c_N such that for any $x \in H$,

$$\sum c_n y_n < \sum c_n \Lambda_n x$$

In particular,

$$\sum c_n y_n < \sum c_n (\Lambda_n \circ f)$$

Now integrating the right hand side gives

$$\sum c_n y_n < \sum c_n y_n$$

which contradicts the fact that y was not in the convex hull. \square

Thus the integral of a continuous function is defined in a great many circumstances. TODO: ESTABLISH PROPERTIES ALA RUDIN.

Example. Convolution of functions can be realized as a vector valued integral. The normal definition

$$(f * g)(x) = \int f(x - y)g(y) dy$$

By removing the value x , can be realized as the function-valued integral

$$f * g = \int f_y g(y) dy$$

If g is integrable, then we can consider this as integration of the function f_y against the finite measure $g dy$, so provided that the map $y \mapsto f_y$ is continuous, then the weak integral exists. The weak convolution is the normal convolution, because if h is a bounded function, then since $|f| * |g|$ is integrable, $(|f| * |g|)h$ is also integrable, so

$$\int |(f * g)(x)| |h(x)| dx \leq \int (|f| * |g|)(x) h(x) dx < \infty$$

and so we may apply Fubini's theorem to conclude that

$$\int \left(\int f(x - y)g(y) dy \right) h(x) dx = \int \left(\int f(x - y)h(x) dx \right) g(y) dy$$

4.1 The Riesz Calculus

4.1.1 Vector-valued integration

We would like to define the integral of arbitrary measurable functions $f : \Omega \rightarrow X$, where (Ω, μ) is a measure space, and X is some Banach space, to be some element $\int f d\mu \in X$. This is not so easy as in the real case, as we have no canonical order through which we can carry summations.

We would like our definition to have properties that ordinary integrals have, namely

$$\Lambda \left(\int f d\mu \right) = \int (\Lambda f) d\mu$$

for every $\Lambda \in X^*$. This is certainly true of other measures, since it is true of finite sums, and the integral can be fort of as some limit of a net of sums.

We shall define a function $f : \Omega \rightarrow X$ to be *measurable* if Λf is measurable for each $\Lambda \in X^*$, where $\Lambda f = \Lambda \circ f$. This is equivalent to the fact that $(\Lambda \circ f)^{-1}(U)$ is measurable for each measurable subset U of the complex plane, and therefore that f is Borel measurable when X has the weak topology. We shall say that a measurable f is *Pettis Integrable* if there is $x \in X$ for which

$$\Lambda x = \int (\Lambda f) d\mu$$

for all $\Lambda \in X^*$. We say x is the integral of f , and denote x by $\int f d\mu$. It is clear that x is unique, for if y also has the property, then $x - y$ is annihilated by all linear maps, hence $x - y = 0$.

Theorem 4.2. *Suppose (Ω, μ) is a compact measure space with a positive Borel measure, $f : \Omega \rightarrow X$ is a continuous map into a Banach space, and $\overline{\text{co}}(f(\Omega))$ is compact in X , then f is integrable, and $\int f d\mu \in \overline{\text{co}}(f(\Omega))$.*

Proof. Assume, without loss of generality, that μ is a probability measure. Let $H = \text{co}(f(Q))$. For a finite set L is a finite set of functionals, let

$$E_L = \left\{ x \in \overline{H} : \Lambda_i(x) = \int (\Lambda_i f) d\mu \right\}$$

Then $E_L \cap E_{L'} = E_{L \cup L'}$, and each E_L is compact, and therefore has the finite intersection property. If we can show that E_L is non-empty for each L , then $\bigcap E_L$ is non-empty, and $\int f d\mu$ exists. Given L , fix an ordering $(\Lambda_1, \dots, \Lambda_n)$, which defines a function from X to \mathbf{R}^n , and consider $m = (m_1, \dots, m_n)$, where

$$m_i = \int (\Lambda_i f) d\mu$$

We shall show that m is in the convex hull of $(L \circ f)(\Omega)$. If m were not in the hull, then, since $\{m\}$ is a convex set, there is a linear functional which separates m and the convex hull of $(L \circ f)(\Omega)$, and since we know the structure of linear functionals on \mathbf{R}^n , there would be $c_1, \dots, c_n \in \mathbf{R}$ for which

$$\sum c_i v_i < \sum c_i m_i$$

whenever $v = (v_1, \dots, v_n)$ is in the hull. But then

$$\sum c_i [(\Lambda_i f)(\omega)] < \sum c_i m_i$$

and integrating,

$$\sum c_i m_i = \sum c_i \int (\Lambda_i f) < \sum c_i m_i$$

Hence m is in the convex hull. Since L is linear, $m = Ly$ for some $y \in H$. At this y , we have

$$\Lambda_i y = m_i = \int (\Lambda_i f) d\mu$$

so $y \in E_L$. Thus $\int f d\mu$ exists, and it lies in $\overline{\text{co}}(f(\Omega))$. \square

This theorem may be generalized to arbitrary complex valued measures by Jordan decomposition.

Theorem 4.3. *If X is a Banach space, $K \subset X$ is compact, and $\overline{H} = \overline{\text{co}}(K)$, then for each $x \in \overline{H}$ there is a regular Borel probability measure μ for which $x = \int \text{id}_K d\mu$.*

Proof. The Riesz representation theorem identifies $C(K)^*$ with the space R of all Borel probability measures on K . Define $\phi : C(K)^* \rightarrow X$ by $\phi(\mu) = \int \text{id}_K d\mu$. The theorem says that $\phi(R) = \overline{H}$. Certainly $K \subset \phi(R)$, for if δ_x is the unit measure at $x \in K$, then $\phi(\delta_x) = x$, for $\Lambda x = \int \Lambda d\delta_x$. Since R is convex, $\phi(R)$ is convex, and we know from the last theorem that it is a subset of \overline{H} . Thus all we need verify is that $\phi(R)$ is closed. This is a consequence of the fact that R is weak * compact in $C(K)^*$, and ϕ is weak * continuous into the weak topology on X , so $\phi(R)$ is weakly closed, and all weakly closed sets are strongly closed.

To see that R is weak * compact, notice that each probability measure has unit variation, and thus lies within the unit ball of $C(K)^*$, which is weak * compact. For $f \in C(K)^*$, let $E_f = \{\mu : \int f d\mu \geq 0\}$. Then each E_f is weak * closed, since the map $f \mapsto \int f d\mu$ is weak * continuous. Hence $\bigcap E_f$ is weak * closed. Similarly, the set $E = \{\mu : \int 1 d\mu = 1\}$ is weak * closed, and $E \cap \bigcap E_f = R$ is therefore weak * closed.

It is sufficient to show ϕ is weak * continuous at the origin. X has a weak subbasis of neighbourhoods at the origin of the form

$$W = \{x \in X : |\Lambda(x)| < r\}$$

for $r > 0$, $\Lambda \in X^*$. Thus we must verify that the set of μ for which $|\int \Lambda d\mu| < r$ is weak * open. But $\Lambda|_K \in C(K)$, so this set is open almost by definition. Thus ϕ is weak * continuous. \square

Theorem 4.4. *If Ω is a compact Hausdorff space with a positive Borel measure μ , X is a Banach space, and $f : \Omega \rightarrow X$ is continuous, then*

$$\left\| \int f d\mu \right\| \leq \int \|f\| d\mu$$

Proof. For each $\Lambda \in X^*$,

$$\left| \Lambda \left(\int f d\mu \right) \right| = \left| \int \Lambda f d\mu \right| \leq \int |\Lambda f| d\mu \leq \|\Lambda\| \int |f| d\mu$$

If Λ is chosen (by the Hahn Banach theorem) such that $\Lambda(\int f d\mu) = \|\int f d\mu\|$, and such that $\|\Lambda\| = 1$, then we obtain the needed inequality. \square

If A is a Banach algebra, then $\int f d\mu$ has an additional property, namely, for each $M \in A$,

$$M \int f d\mu = \int M f d\mu \quad \left(\int f d\mu \right) M = \int f M d\mu$$

This follows since multiplication on the left and right by x is a bounded linear function. If we let $L : N \mapsto MN$ be the left multiplication operator, and fix $\Lambda \in X^*$, then

$$\Lambda \left(M \left(\int f d\mu \right) \right) = (\Lambda \circ L) \left(\int f d\mu \right) = \int (\Lambda \circ L) f d\mu = \int \Lambda(Mf) d\mu$$

the right case is similar.

A *weakly analytic function* is $f : \mathbf{C} \rightarrow A$ for which $\langle \phi, f \rangle$ is analytic for any choice of $\phi \in A^*$. It is clear that all strongly analytic functions are weakly analytic, which follows because

$$\lim_{w \rightarrow 0} \phi \left(\frac{f(z+w) - f(z)}{w} \right) = \phi \left(\lim_{w \rightarrow 0} \frac{f(z+w) - f(z)}{w} \right)$$

For fun, we shall also prove that all weakly analytic functions are strongly analytic.

Theorem 4.5. *Every weakly analytic function $f : D \rightarrow X$ is strongly analytic.*

Proof. Fix $\phi \in X^*$. Consider a particular contour winding counterclockwise around a point w in the domain, which is at least a unit distance away from w at any point on the contour. If $h, k \in \mathbf{C}$ are small enough that $w+h$ and $w+k$ are contained within the contour, then by the Cauchy integral theorem,

$$\begin{aligned} & \phi \left(\frac{1}{h-k} \left[\frac{f(w+h) - f(w)}{h} - \frac{f(w+k) - f(w)}{k} \right] \right) \\ &= \phi \left(\frac{f(w)}{hk} + \frac{f(w+h)}{(h-k)h} - \frac{f(w+k)}{(h-k)k} \right) \\ &= \frac{1}{2\pi i} \int_C \frac{\phi[f(z)] dz}{[z - (w+h)][z - (w+k)][z - w]} \end{aligned}$$

Find δ such that if $\|h\| < \delta$, the distance between any point on C and $w+h$ is greater than $1/2$. Then, if M is the length of C , and K is the supremum of f on C , then

$$\left| \frac{1}{2\pi i} \int_C \frac{\phi[f(z)] dz}{[z - (w+h)][z - (w+k)][z - w]} \right| \leq \frac{4MK}{2\pi} \|\phi\| = \frac{2MK}{\pi} \|\phi\|$$

Applying the Banach Steinhaus theorem (on X^* , viewing elements of X as elements of X^{**}), we conclude that for all h, k sufficiently small, there exists K such that

$$\left| \frac{f(w+h) - f(w)}{h} - \frac{f(w+k) - f(w)}{k} \right| \leq K|h-k|$$

By the completeness of X , the quotients of h and k converge to a well defined quantity as $h-k$ converges to zero. \square

Part II

Linear Analysis

Chapter 5

Topological Vector Spaces

Basic functional analysis deals with norm spaces. But in many problems, more qualitative notions of convergence are encountered that cannot be induced by a norm structure. Here are some examples:

- We say a sequence of functions $\{f_i : X \rightarrow Y\}$ converges *pointwise* to a function f if $f_i(x) \rightarrow f(x)$ for each $x \in X$.
- If X is a locally compact, but non-compact topological space, then we say a sequence of functions $\{f_i : X \rightarrow \mathbf{C}\}$ converges *locally uniformly* to $f : X \rightarrow \mathbf{C}$ if for any compact set $K \subset X$, f_i converges to f uniformly on K .
- Given a norm space X , we say a sequence $\{x_i\}$ in X converges *weakly* to some x if for each linear functional $\phi \in X^*$, $\phi(x_i) \rightarrow \phi(x)$. Similarly, we say $\{\phi_i\}$ in X^* converges in the *weak * topology* to some ϕ if for any $x \in X$, $\phi_i(x) \rightarrow \phi(x)$.

The notions of convergence considered in these examples cannot be induced by a norm, but they can be induced by a more qualitative ‘topological structure’. Our goal in this part of notes is to study vector spaces equipped with this more general topological structure.

In order to be completely general, we begin by making a framework to study function spaces with a topology under the weakest assumptions. A *topological vector space* is a vector space equipped with a topology making both addition and scalar multiplication continuous. In terms of nets, this continuity means exactly that if $x_\alpha \rightarrow x$, $y_\beta \rightarrow y$, and $\lambda_\gamma \rightarrow \lambda$, then

$\lambda_\gamma(x_\alpha + y_\beta) \rightarrow \lambda(x + y)$. In terms of neighbourhoods, the continuity says that if V is an open subset containing $\lambda(x + y)$, then there exists $\varepsilon > 0$, and neighbourhoods W of x and U of y such that if $|\lambda - \gamma| < \varepsilon$, then $\gamma(W + U) \subset V$.

For each $x' \in X$, and $\lambda \neq 0$, the maps $T_{x'}(x) = x + x'$ and $M_\lambda(x) = \lambda x$ are homeomorphisms of X , as a consequence of the continuity of addition and multiplication. As a consequence, every neighbourhood of a point x can be written as $x + W$, where W is a neighbourhood of the origin. Thus most of the topological structure of a space can be understood by looking at the structure around the origin. Another special case of the continuity of addition and multiplication is that, since $1 \cdot (0 + 0) = 0$, for each neighbourhood V of the origin, we can find a subneighbourhood W with $W + W \subset V$. We can repeat this process to find neighbourhoods W such that $W + W + \cdots + W \subset V$ for any fixed, finite number of additions.

Remark. A topological vector space is, of course, an abelian topological group. However, unlike in the case of topological groups, all but the finite dimensional topological vector spaces fail to be locally compact. This is why these two theories proceed down very different paths.

Example. If X is a set, consider the space \mathbf{F}^X of functions from X to \mathbf{F} , under the product topology. Then \mathbf{F}^X is a vector space, which we can make topological if we equip the space with the topology of pointwise convergence. The continuity of the operations follows because if $f_\alpha \rightarrow f$, $g_\beta \rightarrow g$, and $\lambda_\gamma \rightarrow \lambda$, then for each $x \in X$,

$$\lambda_\alpha(f_\alpha(x) + g_\beta(x)) \rightarrow \lambda(f(x) + g(x))$$

because addition and multiplication is continuous on \mathbf{F} , and this shows that the functions $\lambda_\gamma(f_\alpha + g_\beta)$ converges pointwise to $\lambda(f + g)$. The natural neighbourhood base around the origin are the sets

$$\varepsilon B_{x_1, \dots, x_n} = \{f \in \mathbf{F}^X : |f(x_1)|, \dots, |f(x_n)| < \varepsilon\}$$

for $x_1, \dots, x_n \in X$, and $\varepsilon > 0$. This space is only first countable if X is countable.

Example. Consider the space $L_0(X)$ of all measurable functions on a measure space X with the topology induced by convergence of measure. That is, a net f_α converges to f if, for any $\varepsilon > 0$, $\mu(|f_\alpha - f| \geq \varepsilon) \rightarrow 0$. Under a probabilistic interpretation, the probability that the two functions differ by more than ε

converges to zero for every ε . The neighbourhoods of the origin with respect to this topology are of the form

$$\varepsilon U_\delta = \{f \in L_0(X) : \mu(|f| < \varepsilon) \leq \delta\}$$

The space is therefore first countable.

Example. We would expect the space $L_0(X)$ to have an additional topological structure given by pointwise convergence almost everywhere. That is, we would like to have a topology such that $f_\alpha \rightarrow f$ if the set of points x where $f_\alpha(x)$ does not converge to $f(x)$ has zero measure. Tempting a notion as it is, there is no topological structure giving this notion of convergence on $L_0(X)$. To see why, consider the space $L_0[0,1]$, induced by the Lebesgue measure on $[0,1]$. Consider the typewrite sequence of characteristic functions χ_n , defined in terms of the sets $[H_n, H_{n+1}]$ modulo 1, where H_n is the n th Harmonic number. Since $H_n \rightarrow \infty$, every point intersects infinitely many of the intervals, so that χ_n does not converge anywhere to zero. But χ_n does converge in measure to 0, so that for every subsequence of χ_n , there is a further subsequence consisting of functions which converge pointwise almost everywhere to 0. But this contradicts the topological fact that a net x_α converges to a point x if and only if every subnet of x_α has a further subnet which converges to x . Thus there can be no topology whose convergence agrees with almost everywhere convergence.

Example. Consider the space $\mathbf{F}[[X]]$ of formal power series in a single variable. Define a topology on $\mathbf{F}[[X]]$ by letting $f_\alpha \rightarrow f$ if, eventually, the lower order coefficients in the power series expansion of $f_\alpha - f$ vanish. Thus a neighbourhood basis about the origin consists of the decreasing family $\mathfrak{m}, \mathfrak{m}^2$, and so on, where \mathfrak{m} is the ideal consisting of power series with vanishing constant coefficient. It is easy to see that addition is continuous. Multiplication of two power series is also continuous, for if $f_\alpha \rightarrow f$ and $g_\beta \rightarrow g$, then

$$f_\alpha g_\beta - fg = f_\alpha(g_\beta - g) + (f_\alpha - f)g$$

and because each \mathfrak{m}^n is an ideal, if $g_\beta - g$ and $f_\alpha - f$ are in \mathfrak{m}^n , then $f_\alpha g_\beta$ is also in \mathfrak{m}^n . However, scalar multiplication is not continuous, for the convergence properties of this space are ‘too discrete’. Indeed, if we consider $\lambda_n \neq 0$ with $\lambda_n \neq 0$, then $\lambda_n X$ does not converge to 0, because $\lambda_n X \notin \mathfrak{m}^2$. Thus $\mathbf{F}[[X]]$ is a topological ring (over any field), but is not a topological vector space. It is a useful space in algebra, which favours the discrete properties of the topological

structure, but it is not so useful in analysis. We can weaken the topology on $\mathbf{C}[[X]]$ by taking the topology of pointwise convergence of the coefficients, and this gives $\mathbf{C}[[X]]$ a vector space structure, but it no longer reflects the algebraic structure of the space as strongly.

5.1 Basic Theorems

We now introduce some terminology from analytical linear algebra, which will become even more prevalent in the theory of topological vector spaces.

- A *convex* subset C of a vector space X is a set such that if $x, y \in C$ imply $tx + (1 - t)y \in C$ for $0 \leq t \leq 1$, so C contains line segments between each of its points, or in set notation, $tC + (1 - t)C \subset C$ for all $0 \leq t \leq 1$. For any $t, u \in \mathbf{R}$, we have $tC + uC \subset (t + u)C$ which reduces to the ordinary convex equation by dividing both sides by $t + u$.
- A subset B is *balanced* if $\lambda B \subset B$ for $|\lambda| \leq 1$. For any scalar λ , λB is balanced, and if $|\lambda| = 1$, then $\lambda B = B$, since $B = \lambda^{-1}(\lambda B) \subset \lambda B \subset B$ which implies $B = \lambda B$.
- A subset A is *absorbing* if

$$X = \bigcup_{t>0} tA$$

so every $x \in X$ is equal to ty for some scalar $t > 0$ and $y \in A$.

- A set K is (*Von Neumann*) *bounded* if, for every neighbourhood U of the origin, there is $u > 0$ such that if $t > u$, $tU \supset K$.

One of the cornerstones of the basis theory of topological vector spaces is that analytic properties of a space can be established from the existence of open sets with the geometric properties above.

Theorem 5.1. *Every neighbourhood of the origin is absorbing.*

Proof. Given a topological vector space X , and a neighbourhood U of the origin. For any point x , by continuity,

$$\lim_{\lambda \downarrow 0} \lambda x = \left(\lim_{\lambda \downarrow 0} \lambda \right) x = 0 \cdot x = 0$$

Thus, eventually, $\lambda x \in U$, which implies $x \in \lambda^{-1}U$. □

Theorem 5.2. *Every topological vector space has a basis of balanced sets.*

Proof. Given a neighbourhood U of the origin in a topological vector space X , note that by continuity of multiplication there is $\varepsilon > 0$ and $V \subset U$ such that if $|\lambda| < \varepsilon$, $\lambda V \subset U$. Then $\bigcup_{|\lambda| < \varepsilon} \lambda V \subset U$ is a balanced open neighbourhood of the origin contained in U . But this implies that the balanced open neighbourhoods form a basis of the entire topology at the origin. \square

Note that now if $0 < t_1 < t_2 < \dots \rightarrow \infty$, and V is a balanced neighbourhood, then $t_n V$ is balanced, hence contains uV for all $|\lambda| < t_n$ and therefore, since V is absorbing,

$$\bigcup_{t_n} V = \bigcup_{u>0} uV = X$$

The previous theorem says that all balanced neighbourhoods contain absorbing neighbourhoods, so this statements remains true for any neighbourhood V , not just the balanced ones.

Proposition 5.3. *If a space has a bounded open set, it is first countable.*

Proof. Assume the bounded open set V lies at the origin. Then if t_n are a sequence of scalars tending to zero, then $t_n V$ is a neighbourhood basis of the origin, because if W is an arbitrary neighbourhood, then $V \subset t_n^{-1} W$ for some t_n , so $t_n V \subset W$. Thus the $t_n V$ provide a countable neighbourhood base around the origin. \square

Lemma 5.4. *If U is a neighbourhood of the origin such that $x + U$ does not contain y , then there is a neighbourhood V of the origin such that $x + V$ does not intersect $y + V$.*

Proof. Find a balanced subset U_0 of U such that $U_0 + U_0 \subset U$. If there is $u, w \in U_0$ such that $x + u = y + w$, then

$$y = x + u - w \in x + U_0 - U_0 = x + U_0 + U_0 \subset x + U$$

hence $x + U_0$ is disjoint from $y + U_0$. \square

Corollary 5.5. *T_1 topological spaces are Hausdorff.*

From now on, we will assume all the vector spaces we consider are Hausdorff. To see that a particular vector space is Hausdorff, it suffices to show that the intersection of a neighbourhood basis of the origin consists of only the origin, so that $\bar{0} = 0$. In most cases, any topological vector space can be made Hausdorff by considering an appropriate quotient space, ala identifying functions almost everywhere in the theory of L^p spaces.

Corollary 5.6. *If K is a compact set disjoint from C , then there is a neighbourhood W such that $K + W$ is disjoint from $C + W$.*

Proof. For each $x \in K$, we may find a balanced set U_x such that $x + U_x + U_x$ is disjoint from $C + U_x$. By compactness, finitely many U_{x_1}, \dots, U_{x_n} cover K , and then if $W = U_{x_1} \cap \dots \cap U_{x_n}$, we find $K + W$ is disjoint from $C + W$. \square

A corollary of this is that Hausdorff topological spaces are automatically regular. Since $C + W$ is open, we find that, in fact, $\bar{K} + \bar{W}$ does not intersect $C + W$. If we choose $K = \{0\}$, and $C = V^c$ for V a neighbourhood of the origin, we find a neighbourhood W with $\bar{W} \subset V$. Thus for any neighborhood V we can find a subneighbourhood whose closure is contained in V .

Theorem 5.7. *Let X be a topological vector space.*

- (i) *If V is a neighbourhood of the origin, and A is any set, then $\bar{A} \subset A + V$, and in fact $\bar{A} = \bigcap_V A + V$.*
- (ii) *If $A, B \subset X$, then $\bar{A} + \bar{B} \subset \overline{A + B}$.*
- (iii) *The closure of a subspace is a subspace.*
- (iv) *If C is a convex, then C° and \bar{C} are convex subsets.*
- (v) *If B is balanced, then so is \bar{B} , and if $0 \in B^\circ$, then B° is balanced.*
- (vi) *If E is a bounded subset of X , then so is \bar{E} .*
- (vii) *If A and B are bounded, then so is $A + B$.*
- (viii) *If A and B are compact, so is $A + B$.*
- (ix) *If A is compact, and B is closed, $A + B$ is closed.*

Proof. $x \in \overline{A}$ if and only if $x + V$ intersects A for every neighbourhood V of the origin, hence $x \in A - V$, and $-V$ is a neighbourhood of the origin if and only if V is. This proves (i). The continuity of addition essentially proves (ii), for if $x \in \overline{A}$, and $y \in \overline{B}$, then we find $x_\alpha \rightarrow x$, $y_\alpha \rightarrow y$ with $x_\alpha, y_\alpha \in X$, then $x_\alpha + y_\alpha \rightarrow x + y$, and $x_\alpha + y_\alpha \in A + B$. If Y is a subspace, then (ii) implies $\overline{Y} + \overline{Y} \subset \overline{Y + Y} = \overline{Y}$, and $\lambda \overline{Y} = \overline{\lambda Y} = \overline{Y}$, so \overline{Y} is closed under addition and multiplication, and is therefore a subspace. Similarly, if C is convex, then

$$\lambda \overline{C} + (1 - \lambda) \overline{C} = \overline{\lambda C + (1 - \lambda) C} \subset \overline{C}$$

and if $x, y \in C^\circ$, we can find a neighbourhood U such that $x + U, y + U \subset C$, and then for any λ , by convexity

$$(\lambda x + (1 - \lambda)y) + (\lambda U + (1 - \lambda)U) = \lambda(x + U) + (1 - \lambda)(y + U) \subset C$$

which implies $\lambda x + (1 - \lambda)y \in C^\circ$, hence (iv) is shown. If $\lambda B \subset B$ for $|\lambda| < 1$, then $\lambda \overline{B} = \overline{\lambda B} \subset \overline{B}$, so \overline{B} is balanced. Conversely, λB° is open for each $\lambda \neq 0$, so for nonzero $|\lambda| < 1$, $\lambda B^\circ \subset B^\circ$, and provided $0 \in B^\circ$, we find that $\lambda B^\circ \subset B^\circ$ for all $|\lambda| < 1$. If E is a bounded subset, and V is some neighbourhood, then there is some W such that $\overline{W} \subset V$, and if we find t such that $E \subset uW$ for all $u > t$, then $\overline{E} \subset u \overline{W} \subset uV$, hence \overline{E} is bounded. If A and B are bounded, then $A + B$ is bounded, because if W is fixed, we can find U with $U + U \subset W$, and if $A, B \subset \lambda U$, then

$$A + B \subset \lambda(U + U) \subset \lambda W$$

If A and B are compact, then $A + B$ is compact, as it is the image of $A \times B$ under the continuous addition map. If A is compact, and B is closed, then $A + B$ is closed, because if $x_\alpha + y_\alpha \rightarrow z$ in $A + B$, then x_α has a subnet converging to some $x \in A$, and then $y_\alpha \rightarrow z - x$, hence $z - x \in B$, and therefore $z = x + (z - x) \in A + B$. \square

5.2 Convexity and Metrizability

We say a space is *locally convex* if it has a locally convex base.

Theorem 5.8. *Every convex neighbourhood contains a subneighbourhood which is convex and balanced.*

Proof. If U is a convex neighbourhood, consider the set $A = \bigcap_{|\lambda|=1} \lambda U$, which is, convex, and balanced because if $|\gamma| < 1$,

$$\gamma A = \gamma \bigcap_{|\lambda|=1} \lambda U = \bigcap_{|\lambda|=1} \gamma \lambda U = \bigcap_{|\lambda| < |\gamma|} \lambda U \subset A$$

If $V \subset U$ is balanced, then $V \subset A$, so $0 \in A^\circ$, and in particular, A° is a balanced, convex neighbourhood of the origin. \square

Given a set A , we let $\text{Conv}(A)$ denote the *convex hull* of A , the smallest convex set containing A . It can be written as

$$\text{Conv}(A) = \bigcup_{n=1}^{\infty} \bigcup_{t_1 + \dots + t_n = 1} (t_1 A + \dots + t_n A)$$

It follows that the convex hull of an open set is open.

Theorem 5.9. *In a locally convex space, the convex hull of a bounded set is bounded.*

Proof. If W is a convex neighbourhood, and A is bounded, then $A \subset tW$ for some t , and so $\text{Conv}(A) \subset tW$. \square

We now move on to a very important point of the basis theory of topological vector spaces: When can we metrize the topology of a vector space. The most useful metrics on a vector space are the *translation invariant* ones, which satisfy $d(a+x, a+y) = d(x, y)$ for all points a, x, y in the vector space. It turns out that every first countable vector space has such a metrization.

Theorem 5.10. *If X is a first countable vector space with a countable base, then X is metrizable by an invariant metric d whose open balls are balanced. If X is locally convex, the metric can be chosen so the open balls are convex.*

Proof. We consider a balanced neighbourhood base of the origin V_1, V_2, \dots with $V_{n+1} + V_{n+1} + V_{n+1} \subset V_n$. If X is locally convex, choose the V_n to be convex as well. With each Dyadic rational number $r = \sum c_n 2^{-n}$ in $[0, 1]$, we can associate the interval $V_r = \sum c_n V_n$, with $V_1 = X$. These neighbourhoods are monotonic because of the inclusion properties of these neighbourhoods, and satisfy $V_r + V_s \subset V_{r+s}$ (this is most easily proved when $r = s$, in which case the general case is proved). Thus we can define

$$f(x) = \inf\{r \in [0, 1] : x \in V_r\}$$

and $d(x, y) = f(x - y)$. It is clear d is invariant, and satisfies the triangle inequality, because if $x - y \in V_r$ and $y - z \in V_s$, then $x - z \in V_r + V_s \subset V_{r+s}$. Since the V_r is balanced if the V_n are, the open balls of d are balanced, and if the V_r are convex, the open balls of d are convex. \square

Remark. The metric d obtained in the proof above is always bounded, i.e. $0 \leq d \leq 1$, and so we cannot use the same techniques to show that certain topological vector spaces are normable. We will describe the technique to obtain a norm very shortly.

The notion of Cauchy sequences can be used to determine if the metric d is complete. But even if we don't have a metric, we can still define the notion of a Cauchy sequence. We say x_α is a Cauchy net in a topological vector space if, for every neighbourhood U of the origin, eventually we have $x_\alpha - x_\beta \in U$. If x_α converges, then it is obviously Cauchy. Of course, if d is an invariant metric for the topological vector space, these notions agree exactly, which gives the surprising fact that if a single metric for a metrizable topological vector space is complete, then *any* metric for the vector space is complete.

Theorem 5.11. *Cauchy sequences are bounded.*

Proof. Let x_1, x_2, \dots be Cauchy. If U is a balanced neighbourhood of the origin, then eventually for $n \geq N$, $x_n - x_N \in U$. We know that for some λ , $x_1, \dots, x_N \in \lambda U$, hence all of the x_n are in $(\lambda + 1)U$. \square

Theorem 5.12. *A complete metrizable subspace of a topological vector space is closed.*

Proof. If Y is complete and metrizable, and $y_\alpha \rightarrow x$, then y_α is Cauchy, hence converges in Y , and since X is Hausdorff, Y is closed. \square

A *seminorm* on a space X is a positive function $\|\cdot\|$ on X such that $\|x + y\| \leq \|x\| + \|y\|$, and $\|\alpha x\| \leq |\alpha| \|x\|$. These seminorms would give X a metric by setting $d(x, y) = \|x - y\|$, except that we might have $d(x, y) = 0$ with $x \neq y$. One way to fix this is to consider the set Y of all x with $\|x\| = 0$. Then Y forms a subspace of X , and if $y \in Y$, then $\|x + y\| \leq \|x\|$, and also $\|x\| = \|x + y - y\| \leq \|x + y\|$, hence $\|x + Y\| = \|x\|$, and so the seminorm descends to X/Y , and is now a norm. This is the perspective taken in the L^p spaces, where X is the space of all p integrable functions, and Y is the space of p integrable functions equal to zero almost everywhere.

We shall now describe another way to give X a Hausdorff topological structure if, instead of a single seminorm, we have a *family* of seminorms. Given such a family, we can give X a topological structure by letting $x_\alpha \rightarrow 0$ if $\|x_\alpha\| \rightarrow 0$ for every seminorm $\|\cdot\|$ in the family. A family of seminorms is called *separating* if for every $x \neq 0$, there is a seminorm $\|\cdot\|$ with $\|x\| \neq 0$. This implies that the topological structure is Hausdorff, for if x_α converges to two points x and x' , then for any seminorm $\|\cdot\|$

$$\|x - x'\| \leq \|x - x_\alpha\| + \|x_\alpha - x'\| \rightarrow 0$$

hence $\|x - x'\| = 0$, and since this is true for every seminorm, $x = x'$. Every space of this form is locally convex, because a neighbourhood basis is given by the family of convex sets

$$\varepsilon B_{\rho_1 \dots \rho_n} = \{x : \rho_1(x), \dots, \rho_n(x) < \varepsilon\} = \varepsilon(B_{\rho_1} \cap \dots \cap B_{\rho_n})$$

and each of these neighbourhoods is convex, since they are the intersections of convex balls with respect to an individual seminorm. In other words, the family of balls $\varepsilon B_\rho = \{x : \rho(x) < \varepsilon\}$ form a subbase for a topology. If the family of seminorms defining the topology is closed under taking the maximum operation, then these balls actually form a base for the topology.

If $\|\cdot\|$ is a seminorm, then the ‘unit ball’ with respect to this seminorm is convex, balanced, and absorbing. Conversely, suppose A is a convex subset of a topological vector space containing the origin. Then we can define the *Minkowski functional*

$$\mu_A(x) = \inf\{t \geq 0 : x \in tA\}$$

If $x \in \alpha A$ and $y \in \beta A$, then by convexity,

$$\lambda x + \gamma y \in \lambda \alpha A + \gamma \beta A \subset (\lambda \alpha + \gamma \beta) A$$

Taking infima, we find $\mu_A(\lambda x + \gamma y) \leq \lambda \mu_A(x) + \gamma \mu_A(y)$, and so, as a special case, we find $\mu_A(x+y) \leq \mu_A(x) + \mu_A(y)$. If A is an absorbing set, μ_A is always finite, and if, in addition, A is balanced, then $\mu_A(\lambda x) = |\lambda| \mu_A(x)$, so μ_A is a seminorm. If $A_0 = \mu_A^{-1}([0, 1))$ and $A_1 = \mu_A^{-1}([0, 1])$, then $A_0 \subset A \subset A_1$, and $\mu_{A_0} = \mu_A = \mu_{A_1}$. Thus every *continuous* Minkowski functional can be given in terms of an open or a closed set.

The connection between convex sets and Minkowski functionals gives a close correspondence between locally convex spaces and spaces with metrics given by seminorms. If $\|\cdot\|$ is any continuous seminorm on X , we let A consist of the points x with $\|\cdot\| \leq 1$. Then A is convex, because if $\|x\|, \|y\| \leq 1$, then

$$\|\lambda x + (1 - \lambda)y\| \leq \lambda\|x\| + (1 - \lambda)\|y\| \leq \lambda + (1 - \lambda) = 1$$

It is absorbing, since $x \in \|x\|A$, and balanced, since $\|\lambda x\| = |\lambda|\|x\|$. We find

$$\mu_A(x) = \inf\{\lambda \geq 0 : x \in \lambda A\} = \inf\{\lambda \geq 0 : \|x\| \leq \lambda\} = \|x\|$$

so every seminorm is given by a Minkowski functional.

The previous paragraphs hint that the topology of *every* locally convex space is given by a family of seminorms. Indeed, suppose that a space has a convex balanced neighbourhood basis U_α . If $\rho_\alpha = \mu_{U_\alpha}$, then each ρ_α is continuous, because if $x_\beta \rightarrow x$, then eventually $x_\beta - x \in \varepsilon U_\alpha$ for every $\varepsilon > 0$, hence eventually $\rho_\alpha(x_\beta - x) < \varepsilon$. This implies

$$\rho_\alpha(x_\beta) \leq \rho_\alpha(x_\beta - x) + \rho_\alpha(x) \leq \rho_\alpha(x) + \varepsilon$$

$$\rho_\alpha(x_\beta) \geq \rho_\alpha(x) - \rho_\alpha(x - x_\beta) \geq \rho_\alpha(x) - \varepsilon$$

So $\rho_\alpha(x_\beta) \rightarrow \rho_\alpha(x)$. Conversely, if $\rho_\alpha(x_\beta) \rightarrow 0$ for every fixed α , then eventually $\rho_\alpha(x_\beta) \leq 1$, so $x_\beta \in U_\alpha$, so x_β is eventually in U_α for every α , so $x_\beta \rightarrow 0$. This shows that the topology of the vector space is exactly the same as the one formed from the family of seminorms ρ_α .

If the topology of a space X is specified by a countable family of seminorms, then the space is first countable, and therefore metrizable. However, in this situation we can come up with a natural metric to use, rather than appealing to the theorem we just came up with. If ρ_1, ρ_2, \dots is an ordering of the seminorms, then the functions

$$d_n(x, y) = \frac{\rho_n(x - y)}{1 + \rho_n(x - y)}$$

are bounded and satisfy the triangle inequality, and the metric for the space can be given by

$$d(x, y) = \sum_{n=1}^{\infty} \frac{d_n(x, y)}{2^n} = \sum_{n=1}^{\infty} \frac{\rho_n(x - y)}{2^n(1 + \rho_n(x - y))}$$

Since *any* invariant metric on a complete, metrizable space is complete, there is normally not a need to switch to a different metric when doing general calculations on metric spaces, unless it is possible to switch to a single norm for the entire space.

Theorem 5.13. *Let X be a locally convex, specified by a family of seminorms ρ_α . Then a set E is bounded if and only if ρ_α is bounded on E for each α .*

Proof. Suppose E is bounded. If U_α is the unit ball with respect to ρ_α , then for some $t > 0$, $E \in tU_\alpha$, so $\rho_\alpha(E) \leq t$. Conversely, suppose that $\rho_\alpha(E) \leq t_\alpha$ for each α . Then we consider a neighbourhood basis of the elements $\varepsilon U_{\rho_1 \dots \rho_N}$, as introduced above, and for each such element we find

$$E \in \max(t_1, \dots, t_N) \varepsilon U_{\rho_1 \dots \rho_N}$$

so E is bounded. □

To finish our understanding of locally convex spaces, we ask when the topology of a locally convex space can be given by a single norm. If $\|\cdot\|$ is a norm for the space, then the unit ball with respect to this norm is bounded by the theorem above, and convex. It turns out this is sufficient for a locally convex space to be normable.

Theorem 5.14. *A space is normable if and only if it has a bounded, convex neighbourhood.*

Proof. If U is a bounded, convex neighbourhood about the origin, then the εU form a neighbourhood basis around the origin. It follows that the topology of the space is specified by the single Minkowski functional μ_U . The fact that μ_U is actually a norm, rather than a seminorm, follows from the fact that the space is Hausdorff. □

Example. Consider the class $C(X)$ of continuous functions on a locally compact space X , under the topology of locally uniform convergence (uniform convergence on compact sets). Then $C(X)$ is a topological vector space. For each compact set K , we have the seminorms $\|\cdot\|_{L^\infty(K)}$, which define a seminorms which show $C(X)$ and $H(X)$ have the topology of a locally convex space. If X is σ compact, then the topology can be given by a countable family of seminorms, and so $C(X)$ is metrizable. If we write $X = \lim_{n \rightarrow \infty} K_n$, and f_n is Cauchy, then the f_n converge uniformly to some continuous function g_m on K_m , and

g_{m+1} agrees with g_m on K_m , hence we can consider the function g defined everywhere by letting $g(x) = g_m(x)$ for all $x \in K_m$, and then $f_n \rightarrow g$. Thus $C(X)$ is Fréchet. However, $C(X)$ is not normable if X is not compact. We know a set $E \subset C(X)$ is bounded if and only if there is M_K such that $\|f\|_{L^\infty(K)} \leq M_K$ for each $f \in E$. But we have a neighbourhood consisting of $U_{K_1 \dots K_N \varepsilon_1 \dots \varepsilon_N}$. And since X is not compact, there is some $x \notin K_1 \cup \dots \cup K_N$, and since X is completely regular, we can find a continuous function f with $f(x) = 1$, but f vanishing on $K_1 \cup \dots \cup K_N$. It follows that $nf \in U_{K_1 \dots K_N \varepsilon_1 \dots \varepsilon_N}$ for all integers n , and the nf are not bounded. Thus $C(X)$ is not normable unless X is compact.

Example. If X is an open subset of \mathbf{C} , then the space $H(X)$ of holomorphic functions on X is also a topological vector space under locally uniform convergence. Since the local uniform limit of a holomorphic function is a holomorphic function, $H(X)$ is a closed subspace of $C(X)$, and is therefore a Fréchet space. Montel's theorem in the theory of complex variables says that $H(X)$ has the Heine-Borel property: Every closed, bounded subset of $H(X)$ is compact. But since $H(X)$ is never finite dimensional, $H(X)$ cannot be locally compact, hence $H(X)$ cannot be locally bounded. It follows that $H(X)$ is also not normable.

Example. Let X be an open subset of \mathbf{R}^n . Then we can give $C^\infty(X)$ the structure of a topological vector space by letting $f_\alpha \rightarrow f$ if all the partial derivatives of f_α converge to the partial derivatives of f locally uniformly. Arguing essentially the same as with $C(X)$ and $H(X)$, this shows $C^\infty(X)$ is a metrizable locally convex space. If f_α is a Cauchy sequence in this space, then f_α is Cauchy on each compact set K , hence f_α and all its partial derivatives converges locally uniformly on K to some C^∞ function g_K . Putting these things together gives that $C^\infty(X)$ is a Fréchet space. Surprisingly $C^\infty(X)$ also satisfies the Heine-Borel property. This is easily proven using the Arzela-Ascoli theorem. Since $C^\infty(X)$ is never finite dimensional, this space cannot be locally bounded, hence $C^\infty(X)$ isn't normable.

5.3 Linear Maps

As in the functional analysis of Banach spaces, linear maps play an important role in the theory. In particular, we are interested in those linear maps between topological vector spaces which are continuous. Continuous maps are exactly those which are uniformly continuous, because if $T : X \rightarrow Y$ is continuous, then for every neighbourhood V of the origin in

Y , there is a neighbourhood U of the origin in X such that if $x - y \in U$, then $Tx - Ty \in V$.

Lemma 5.15. *If T is continuous at 0, then T is continuous everywhere.*

Proof. Suppose T is continuous at the origin. If $x_\alpha \rightarrow x$, then $x_\alpha - x \rightarrow 0$, hence $T(x_\alpha - x) = T(x_\alpha) - T(x) \rightarrow 0$, so $T(x_\alpha) \rightarrow T(x)$. Thus T is continuous everywhere. \square

The theory of linear functions is difficult to classify in the general case, but linear *functionals* are easier to understand. Given X , we let $X^\#$ denote the class of linear functionals, and X^* the class of *continuous* linear functionals.

Theorem 5.16. *If $f \in X^\#$ is a non-zero linear functional, then the following are equivalent:*

1. f is continuous.
2. The nullspace of f is closed.
3. The nullspace of f is not dense in X .
4. f is bounded in some neighbourhood of the origin.

Proof. (1) implies (2) implies (3) is obvious. If (3) holds, so the nullspace of f is not dense in X , we can find $x \in X$ and a balanced set U such that $0 \notin f(x + U) = f(x) + f(U)$, hence $f(u) \neq -f(x)$ for any $u \in U$. Since U is balanced, this must imply that $f(u) \neq \lambda f(x)$ for any $|\lambda| \geq 1$, so $|f(u)| \leq |f(x)|$ for all $u \in U$, so f is bounded on U , implying (4). If f is bounded in some neighbourhood of the origin, then $|f(u)| \leq M$ for all $u \in U$, and if $x_\alpha \rightarrow 0$, then eventually for all $\varepsilon > 0$, $x_\alpha \in \varepsilon U$, hence $|f(x_\alpha)| \leq \varepsilon M$, implying $f(x_\alpha) \rightarrow 0$. This means f is continuous at the origin. The general statement now follows from the fact that a linear function is continuous if and only if it is continuous at the origin. \square

Example. Suppose X is a locally convex space, with the topology specified by a family of seminorms closed under the maximum operation. Then the balls εU_ρ form a basis for the space. It follows that if $f \in X^*$ is a continuous linear functional, then f is bounded on a neighbourhood on the origin, so there exists

ρ, ε , and M such that if $\rho(x) \leq \varepsilon$, then $|f(x)| \leq M$. It then follows that for any $x \in X$,

$$|f(x)| = \frac{\rho(x)}{\varepsilon} \left| f\left(\frac{\varepsilon x}{\rho(x)}\right) \right| \leq (M/\varepsilon)\rho(x)$$

so the functional is directly bounded above by a seminorm.

For finite dimensional topological vector spaces, the continuous linear maps are particularly easy to consider.

Lemma 5.17. *Every linear map $T : \mathbf{F}^n \rightarrow X$ is continuous.*

Proof. If we write $x_i = T e_i$, where e_i is the i 'th element of canonical basis of \mathbf{F}^n , then $T(a) = \sum a_i x_i$, and the continuity follows because addition and scalar multiplication is continuous. \square

Corollary 5.18. *Every finite dimensional vector space has a unique topology making it into a topological vector space.*

Proof. if τ_0 and τ_1 are two topologies on \mathbf{F}^n turning it into a topological vector space, then the identity map on \mathbf{F}^n is a continuous map from (\mathbf{F}^n, τ_0) to (\mathbf{F}^n, τ_1) , and from (\mathbf{F}^n, τ_1) to (\mathbf{F}^n, τ_0) , hence we conclude the identity map is a homeomorphism. \square

Corollary 5.19. *Every finite dimensional subspace of a Hausdorff topological vector space is closed.*

Proof. If $Y \subset X$ is a subspace of X , let $T : \mathbf{F}^n \rightarrow Y$ be an isomorphism. Then we know $S = T(S_{\mathbf{F}^n})$ is compact in X , and $B = T(B_{\mathbf{F}^n})$ is open. By compactness, there is an open neighbourhood V of the origin in X which does not intersect S . Then $W = T^{-1}(V \cap Y)$ is an open neighbourhood of the origin in \mathbf{F}^n disjoint from $S_{\mathbf{F}^n}$, so that it only contains elements in B , and so \overline{W} is compact, hence $V \cap Y \subset T(\overline{W}) \subset Y$, and so the closure of each $V \cap Y$ lies in Y . Now if $z \in \overline{Y}$ is arbitrary, then $z = tx$ for $x \in V$, and since \overline{Y} is a subspace we find x lies in \overline{Y} . But if $y_\alpha \rightarrow x$, with $y_\alpha \in Y$, then y_α is eventually in $Y \cap V$, which is a closed subset of Y , hence we find $x \in Y$, hence $z \in Y$. \square

Theorem 5.20. *Every locally compact vector space X has finite dimension.*

Proof. If a vector space has a precompact neighbourhood V , then that neighbourhood is bounded, so $V/2^n$ form a local base for X . Now there are $x_1, \dots, x_n \in V$ such that

$$\overline{V} \subset (x_1 + V/2) \cup \dots \cup (x_n + V/2)$$

Then the span of the x_i is a closed subspace Y of X , and $\overline{V} \subset Y + V/2$. Then, of course, $V \subset Y + V/2$, hence, dividing by two, we find $V/2 \subset Y + V/4$, and so $\overline{V} \subset Y + V/4$. Continuing this process, we find $\overline{V} \subset \bigcap (Y + V/2^n)$, and since $V/2^n$ form a local base, we find $\overline{V} \subset \overline{Y} = Y$. We conclude that Y contains an absorbing set for X , hence $Y = X$. \square

Since we have a definition of boundedness in any topological vector space X , we can ask if a version of the Heine Borel theorem holds in X . That is, are closed, bounded subsets of X compact.

Corollary 5.21. *No infinite dimensional space can be locally bounded and satisfy the Heine-Borel theorem.*

Proof. If V is a bounded neighbourhood, then it is precompact, hence the space is locally compact, so finite dimensional. \square

In the context of norm spaces, we found a linear map was continuous if and only if it was bounded. An analogous definition of a bounded operator $T : X \rightarrow Y$ between topological vector spaces is an operator that maps bounded sets to bounded sets. It is easy to prove that a continuous map is bounded, but the converse is not always true. However, in a metrizable space, it is true.

Lemma 5.22. *if X is a metrizable space, and $x_n \rightarrow 0$, then there are constants $\lambda_n \rightarrow \infty$ with $\lambda_n x_n \rightarrow 0$.*

Proof. If d is translation invariant, it is easy to prove that for any x , $d(nx, 0) \leq nd(x, 0)$. If $d(x_n, 0) \rightarrow 0$, there are an increasing sequence of positive integers N_m with $d(x_n, 0) \leq 1/m^2$ for $n \geq N_m$. Then $d(mx_{n_m}, 0) \leq 1/m$, so if λ_n is the largest m with $n_m \leq n$, then $\lambda_n x_{n_m} \rightarrow 0$. \square

Example. *If X is non-metrizable, this need not be true. Consider the space $\mathbb{C}^{[0,1]}$ of all complex functions on $[0, 1]$, under the topology of pointwise convergence. Since $[0, 1]$ has the same cardinality as the space $\mathbb{C}^{\mathbb{N}}$ of all complex-valued sequences converging to zero pointwise, we can define a map $F : [0, 1] \rightarrow$*

$\mathbb{C}^{\mathbb{N}}$, and then define a sequence f_n of functions converging pointwise to zero by letting $f_n(x) = F(x)_n$. For any sequence of constants $\lambda_1, \lambda_2, \dots$ converging to ∞ , the sequence $1/\lambda_1, 1/\lambda_2, \dots$ converges to zero, hence there is some $x \in [0, 1]$ with $F(x)_n = 1/\lambda_n$. Then $\lambda_n f_n(x) = \lambda_n/\lambda_n = 1$, so $\lambda_n f_n$ does not converge pointwise to zero.

Theorem 5.23. *If X is a metrizable space, every bounded operator with X as a domain is continuous.*

Proof. Since X is first countable, it suffices to prove that if x_n is a sequence converging to zero, and $T : X \rightarrow Y$ is bounded, then T is continuous. Certainly x_n is bounded, so Tx_n is bounded. If we pick λ_n as in the last lemma, so $\lambda_n x_n \rightarrow 0$, then $\lambda_n x_n$ is bounded, so $\lambda_n Tx_n$ is bounded, and since $\lambda_n^{-1} \rightarrow 0$, $Tx_n = \lambda_n^{-1}(\lambda_n Tx_n) \rightarrow 0$. \square

Example. A quasinorm space is a vector space X equipped with a non-negative real-valued function $\|\cdot\|$ such that $\|x\| = 0$ if and only if $x = 0$, $\|\lambda x\| = |\lambda|\|x\|$, and $\|x + y\| \lesssim \|x\| + \|y\|$, rather than the stronger inequality $\|x + y\| \leq \|x\| + \|y\|$. The most precient examples of quasinorm spaces are the $L^p(X)$ spaces over some measure μ , where $p < 1$. Since the unit ball in these spaces are not convex, we know that $\|\cdot\|_p$ is not a norm. However, the norms corresponding to these spaces are quasinorms. To see this, we use the fairly trivial inequality $(a+b)^p \leq 2^p(a^p + b^p)$, for $a, b, p > 0$, which follows because

$$(a + b)^p \leq (2 \max(a, b))^p = 2^p \max(a, b)^p \leq 2^p(a^p + b^p)$$

so that $\|f + g\|_p^p \leq 2^p(\|f\|_p^p + \|g\|_p^p)$, and therefore

$$\|f + g\|_p \leq 2(\|f\|_p^p + \|g\|_p^p)^{1/p} \leq 2^{1+1/p}(\|f\|_p + \|g\|_p)$$

So the ' L^p norm' is a quasinorm. This quasinorm gives $L^p(X)$ the structure of a complete metric space, the proof of completeness being exactly the same as in the case where $p \geq 1$. The fact that only a single quasinorm is defined in the definition of the topology shows the space is locally bounded. However, the space is not locally convex, hence not a Fréchet space. In fact, as an example, $L^p[0, 1]$ contains no convex open sets other than \emptyset and $L^p[0, 1]$. To see this, let W be an non-empty, open, convex set, and without loss of generality, let $0 \in W$, and suppose W contains all functions f with $\|f\|_p < \varepsilon$. Given any $f \in L^p[0, 1]$, and any integer N , we can divide $[0, 1]$ into disjoint intervals I_1, I_2, \dots, I_N with

$$\|f\|_{L^p(I_n)} = \frac{\|f\|_{L^p[0,1]}}{N^{1/p}}$$

If $f_n = N\chi_{I_n}f$, then $\|f_n\|_{L^p[0,1]} = N^{1-1/p}\|f\|_{L^p[0,1]}$. Since $p < 1$, we can choose N large enough that $\|f_n\|_{L^p[0,1]} < \varepsilon$, hence $f_n \in W$, and then by convexity,

$$f = \frac{f_1 + \cdots + f_N}{N} \in W$$

Thus $W = L^p[0,1]$. Now a strange result of this is that if $T : L^p[0,1] \rightarrow X$ is a continuous linear transformation into a locally convex space X , then $T = 0$. This follows because if W is an open, convex neighbourhood, then $T^{-1}(W)$ is open and convex, hence $T^{-1}(W) = L^p[0,1]$. But then taking W to be an arbitrary convex neighbourhood of the origin,

$$T^{-1}(0) = T^{-1}\left(\bigcap W\right) = \bigcap T^{-1}(W) = L^p[0,1]$$

Thus there are no nonzero continuous linear functions on $L^p[0,1]$.

For non-metrizable locally convex spaces, we still have a nice analytical criterion for continuity in terms of seminorms.

Theorem 5.24. *Let $T : X \rightarrow Y$ be a map between convex spaces. Then T is continuous if and only if for any continuous seminorm $\|\cdot\|_Y$ on Y , there is a continuous seminorm $\|\cdot\|_X$ on X such that $\|Tx\|_Y \leq \|x\|_X$.*

Proof. Suppose T is continuous. Then the unit ball B in the Y norm is open, so $T^{-1}(B)$ is an open neighbourhood of the origin in X , and there therefore contains a convex, balanced subset about the origin which is the unit ball with respect to some continuous seminorm $\|\cdot\|_X$. Thus $\|Tx\|_Y \leq 1$ if $\|x\|_X \leq 1$, which implies that

$$\|Tx\|_Y = \|x\|_X \|T(x/\|x\|_X)\|_Y \leq \|x\|_X$$

Conversely, if the other condition is satisfied and U is an open convex set about the origin in Y , then it is the unit ball with respect to some norm $\|\cdot\|_Y$, and so there is $\|\cdot\|_X$ such that $\|Tx\|_Y \leq \|x\|_X$, hence the unit ball with respect to the X norm is an open neighbourhood which maps into U , so T is continuous at the origin, hence continuous everywhere. \square

If X and Y are generated by a family of seminorms closed under maxima, we can restrict the conditions of the theorem to only the seminorms contained in this group, if we weaken the condition $\|Tx\|_Y \leq \|x\|_Y$ to $\|Tx\|_Y \lesssim \|x\|_X$.

5.4 Quotient Spaces

If Y is a subspace of a topological vector space X , then X/Y can be given the quotient topology, and one easily verifies that addition and multiplication is continuous, so the quotient of a topological vector space is a topological vector space. Since the neighbourhood of $0 + Y$ in X/Y is an open set containing Y ,

$$\overline{0 + Y} = \bigcap_{U+Y} (0 + Y) + (U + Y) = \bigcap_U (U + Y) = \overline{Y} + Y$$

So X/Y is Hausdorff if and only if Y is closed.

We can now argue why assuming a space is Hausdorff isn't too strong a statement. Let X be a (not necessarily Hausdorff) topological vector space. Then 0 is a subspace of X , and so the closure of 0 is also a closed subspace. But now if we consider the space $X/\overline{0}$, we obtain a Hausdorff topological space. Thus the techniques of Hausdorff topological spaces can normally be used to obtain results about non Hausdorff spaces by means of quotienting.

Theorem 5.25. *The quotient map $\pi : X \rightarrow X/Y$ is open, and a neighbourhood base for X/Y is given by taking $\pi(U)$, where U ranges over a neighbourhood base.*

Proof. The fact that $\pi(U)$ is open when U is open follows because

$$\pi^{-1}(\pi(U)) = U + Y$$

is obviously a union of open neighbourhoods. If V is an open neighbourhood of the origin in X/Y , then $\pi^{-1}(V)$ is a neighbourhood of the origin, hence contains some U in a neighbourhood base, and then $\pi(U) \subset V$. \square

Theorem 5.26. *The quotient of a locally convex / locally bounded / metrizable / normable / complete metrizable space also has these properties.*

Proof. The image of a convex / balanced / absorbing / bounded set under a quotient map is also convex / balanced / absorbing / bounded. This shows that the quotient of a locally convex / bounded space is locally convex / bounded. Since the quotient of a first countable space is first countable, this shows that the quotient of a metrizable space is metrizable. The

fact that the quotient of a normable space is normable follows because a normable is precisely a locally convex and locally bounded space. The only remaining fact is to show that the quotient of a complete metrizable space is complete. If d is a complete invariant metric for the space, with $f(x) = d(0, x)$, then the metric for the quotient is given by

$$d(x_0 + Y, x_1 + Y) = \inf\{f((x_1 - x_0) + y) : y \in Y\}$$

The triangle inequality follows because if $f((x_1 - x_0) + y_0) \leq \alpha$ and $f((x_2 - x_1) + y_1) \leq \beta$, then $d(x_1 + y_0, x_0) \leq \alpha$, and $d(x_2 + y_1 + y_0, x_1 + y_0) \leq \beta$, and so the triangle inequality on X gives

$$d(x_0, x_2 + y_1 + y_0) \leq \alpha + \beta$$

so $d(x_0 + Y, x_2 + Y) \leq \alpha + \beta$, and taking infima gives the triangle inequality. If $x_\alpha + Y \rightarrow 0$ in X/Y , then we can pick y_α with $x_\alpha + y_\alpha \rightarrow 0$ in X , hence $d(0, x_\alpha + y_\alpha) \rightarrow 0$, so $d(0, x_\alpha + Y) \rightarrow 0$. Conversely, if $d(0, x_\alpha + Y) \rightarrow 0$, then there are y_α with $d(0, x_\alpha + y_\alpha) \rightarrow 0$, so $x_\alpha + y_\alpha \rightarrow 0$, hence $x_\alpha + Y \rightarrow 0$. Thus d is a metric for the topological structure of X/Y . Now if d is complete on X , it remains to show d is complete on X/Y . Suppose $x_n + Y$ is Cauchy. Then we can inductively, for each n , pick k_n and y_n such that $d(x_{k_{n+1}} + y_{n+1}, x_{k_n} + y_n) \leq 1/2^n$. This means $x_{k_n} + y_n$ is Cauchy, hence converges, hence $x_{k_n} + Y$ converges, and because of the Cauchy property it follows that $x_n + Y$ converges to the same value. \square

Corollary 5.27. *If Y is a closed subspace of X , and Z is finite dimensional, then $Y + Z$ is closed.*

Proof. Z/Y is finite dimension in X/Y , hence is closed. But $Y + Z$ is the inverse image of Z/Y under the quotient map, which is continuous, so $Y + Z$ is closed. \square

5.5 Inductive Limits

We have a theory of topological vector spaces, and maps between them, so we have a category. Let us study inductive limits in this space. The canonical example of this is the test function space $\mathcal{D}(\Omega)$ on an open subset Ω of \mathbf{R}^d , which is the inductive limit of the spaces $C_c^\infty(K)$ of smooth

functions with support on a compact set $K \subset \Omega$. For general (not necessarily Hausdorff) topological spaces, inductive limits exist. Given a system $\{f_{ij} : X_i \rightarrow X_j\}$ of linear maps, we can form the topological inductive limit X , together with the family of maps $f_i : X_i \rightarrow X$ such that $U \subset X$ is open if and only if $f_i^{-1}(U)$ is open in X_i for each i . The set X is defined such that $X = \bigcup f_i(X_i)$, and $f_i(x) = f_j(x)$ if and only if there exists $k > i, j$ such that $f_{ik}(x) = f_{jk}(y)$. This is precisely the topology with the most open sets, such that each of the maps f_i is continuous. The algebraic operations on each of the spaces X_i induces algebraic operations on X , and it is simple to see that these operations are continuous. Thus X is a topological vector space. There are several problems with this construction that make this not so nice, however, because topological properties do not play quite so nicely with direct limits (the direct limit of Hausdorff topological vector spaces, for instance, need not be Hausdorff).

Our construction can be made a bit nicer if we work in the subcategory of *locally convex topological spaces*. Working in this subcategory, we require less functions to be continuous, and so the direct limit can be defined with less open sets. In particular, if each of the spaces $\{X_i\}$ is locally convex, then to form an direct limit in this subcategory, then we obtain an direct limit in the category of locally convex topological vector spaces by declaring a *convex* set $U \subset X$ to be an open neighborhood of the origin if and only if $f_i^{-1}(U)$ is open for each i , and then using these sets as a basis for the overall topology. It is easy to see this is a basis, because the intersection of convex sets is convex, and thus we get a topology. Even in this subcategory, the direct limit of Hausdorff locally convex spaces need not be Hausdorff.

Lemma 5.28. *Consider a countable direct system $X_1 \rightarrow X_2 \rightarrow \dots$, where each X_i is a locally convex T1 Hausdorff space, and such that each map $f_i : X_i \rightarrow X_{i+1}$ is injective. Then the direct limit X is T1 (and thus Hausdorff).*

Proof. Consider $z, w \in X$. Suppose i is the smallest index such that $z, w \in X_i$ and set $U_{i-1} = \emptyset$. Starting with $k = i$, define a convex set U_k , which is open in X_k , such that $U_{k-1} \subset U_k$, $z \in U_k$, but w is not an element of $\overline{U_k}$. To do this, we start by finding a convex, balanced set V such that $w + V$ is disjoint from $\overline{U_{k-1}}$ by assumption. Then $U_k = U_{k-1} + V$ is convex, open, contains z , and its closure is disjoint from w . Then the set $\bigcup_k U_k$ is a convex, open neighborhood of z in X , not containing w , so X is T1. \square

A similar method to the proof above, combined with the following lemma, shows that if each of the maps $X_i \rightarrow X_{i+1}$ is an embedding, then the maps $X_i \rightarrow X$ will also be embeddings. A space obtained from the limit of countably many Fréchet spaces embedded in one another is called an *LF Space*.

Lemma 5.29. *Suppose X and Y are locally convex vector spaces, with Y a subspace of X . Then for any convex open set $V \subset Y$, there exists a convex open set $U \subset X$ such that $U \cap Y = V$.*

Proof. By the definition of the subspace topology, we can write $V = V' \cap Y$ for some open set V' in X . Let U be the convex hull of V' . Because V is convex, $U \cap Y = V$. But U is also open in X . \square

A direct limit of a countable family of Fréchet spaces as in the Lemma above is called an *LF space*. Most interesting results about direct limits are reserved for LF spaces, like the following result. For instance, that every LF space is complete.

These notes are probably worth looking into if I ever want to learn more about limits of convex spaces: Bierstedt, Klaus-Dieter (1988). *An Introduction to Locally Convex Inductive Limits. Functional Analysis and Applications.* Singapore-New Jersey-Hong Kong: Universitätsbibliothek. pp. 35–133. MR 0046004. Retrieved 20 September 2020.

Chapter 6

Weak Topologies

Example. In Banach theory, the most important examples of topological spaces are those obtained by placing other topologies on a norm space than the original norm topology. Given a norm space X , we have a space X^* consisting of continuous linear functionals on X . We say that a net x_α converges weakly to x if $f(x_\alpha) \rightarrow f(x)$ for every $f \in X^*$. This defines the weak topology on X . Similarly, on X^* we may define the weak $*$ topology, by letting $f_\alpha \rightarrow f$ if $f_\alpha(x) \rightarrow f(x)$ for every $x \in X$. Note that we may also put a weak topology on X^* , since it is also a norm space, but the weak and weak $*$ topologies on X^* rarely correspond (the weak topology has more stringent notions of continuity). These topologies are effectively topologies induced by pointwise convergence, because the weak topology is obtained by embedding X in \mathbb{F}^{X^*} , and the weak $*$ topology is obtained by considering X^* as a subspace of \mathbb{F}^X . The weak topology on X is rarely equal to the original topology on X , except in degenerate circumstances. For instance, if we consider $e_n \in l_p$, for $1 < p < \infty$, then $\|e_n - e_m\| = 1$, so e_n does not converge, yet if we consider any element $f \in l_q$ of the dual space, then $\langle f, e_n \rangle = f(n)$, which converges to zero as $n \rightarrow \infty$. Thus e_n converges to zero in the weak topology. In l_1 , e_n does not converge at all, because the function $f(n) = (-1)^n$ is in l_∞ , and $\langle f, e_n \rangle = (-1)^n$ does not converge. It is important to note that the weak $*$ topology depends on the predual we are using. If we consider $l_1 = c_0^*$, then the e_n converges to zero in the weak $*$ topology, yet if we consider l_1 as c^* , then the e_n do not converge in the induced weak $*$ topology.

The weak topology on a space is rarely first countable. If $A = \{\sqrt{n}e_n : n \in \mathbf{N}\}$ is viewed as a subset of l_2 , then 0 is contained in the weak closure of A , because if $\varepsilon > 0$ and we are given a finite set of sequences $\{a_n^k\}_{n \in \mathbf{N}}$, for $1 \leq k \leq K$

and, for any n , there is k_n such that $|a_n^{k_n}|\sqrt{n} > \varepsilon$, and if we consider the sum $\sum_k |a_n^k| \in l^2$, then

$$\sum_n \left| \sum_k |a_n^k| \right|^2 \geq \sum_n |a_n^{k_n}|^2 \geq \varepsilon \sum_n \frac{1}{n} = \infty$$

so there must be an element of A in each $U_{\varepsilon, \{a_n^1\}, \dots, \{a_n^{k_n}\}}$, and hence 0 is in the weak closure of A . However, 0 is not the weak limit of any sequence of elements in A , for

$$\sum_{n=0}^{\infty} \frac{e_n}{n} \in l_2$$

and if we consider any sequence $\sqrt{n_i}e_{n_i}$ that converges weakly to zero, then we find that $(n_i)^{-1/2}$ converges to zero, hence $n_i \rightarrow \infty$. But then the sequence is not norm bounded, and hence cannot converge!

Part III

Operator Theory

Chapter 7

Banach Algebras

In this chapter, we develop the machinery to study the various classes of bounded operators which occur when classifying transformations on a space. To do this, we employ tools from abstract algebra. Recall that an algebra over a field F is a (not necessarily unital) ring A together with a fixed embedding of F into A , which gives A a vector space structure. A consistent norm attached to an algebra is a nice situation to study operators over a Banach space. What's more, the general theory gives light to many other circumstance which would not have been apparent were we just analyzing concrete operators. We shall use capital letters, like M and N , to denote abstract elements of these algebras, and we denote algebras by capital letters near the beginning of the alphabet, such as A or B . A *Banach Algebra* is a Banach space A which is also an algebra, and satisfies $\|MN\| \leq \|M\|\|N\|$ for all $M, N \in A$. A *unital Banach Algebra* is a Banach algebra with an identity $1 \neq 0$, which we will assume also satisfies $\|1\| = 1$. Thus multiplication is a continuous operation, and the norm is normalized to give the multiplicative identity unit norm. We shall pretty much only consider Banach algebras over the complex numbers, because here we can apply the unique properties of complex analysis.

Example. Let X be a locally compact topological space. Then the space $C_b(X)$ of continuous, bounded functions on X forms a unital commutative Banach algebra under the pointwise product, and equipped with the L^∞ norm. The space $C_0(X)$ of continuous functions vanishing at infinity forms a subalgebra of $C_b(X)$, and is unital precisely when X is compact, in which case $C_0(X) = C_b(X) = C(X)$. Regardless, we can add an identity by consider the family of

all functions of the form

$$C_0(X)^\# = \{\lambda + f : \lambda \in \mathbf{C}, f \in C_0(X)\},$$

which can be viewed as the space of all continuous functions which converge to a common value at ∞ , which can thus be naturally identified with $C(X \cup \{\infty\})$, where $X \cup \{\infty\}$ is the one point compactification of X .

As a simple example, if $X = \{x_1, \dots, x_n\}$ is a discrete topological space, then $C(X) \cong \mathbf{C}^n$, where the right hand side is a Banach algebra if we equip it with the standard norm, and the pointwise product, since the Cauchy-Schwartz inequality then tells us that $\|xy\| \leq \|x\|\|y\|$ for all $x, y \in \mathbf{C}^n$.

Example. The space $L^\infty(X)$ of essentially bounded functions on a measure space X is a Banach algebra, and it is unital except in the trivial case where $|X| = 0$, in which case $L^\infty(X) = 0$.

Example. If Ω is a compact set in \mathbf{C} , then we define $A(\Omega)$ to be the subset of $C(\Omega)$ consisting of all functions which are analytic on the interior of Ω . Since uniform convergence preserves holomorphicity, $A(\Omega)$ is a closed subalgebra of $C(\Omega)$, and is therefore a Banach algebra. $A(\mathbf{D})$ is known as the disk algebra.

Example. Let X be a Banach space. Then X has the structure of a somewhat trivial, commutative Banach algebra if we define $x_1 x_2 = 0$ for any $x_1, x_2 \in X$. This is often a useful family of spaces to construct counterexamples, or to show why certain assumptions are necessary in Banach algebra theory.

A Banach algebra A is generated by a set $S \subset A$ if there is no proper closed subalgebra of A containing S . In this language, the Stone-Weierstrass theorem says that, for any locally compact Hausdorff space X , a family of functions $S \subset C_0(X)$ generates $C_0(X)$ if and only if S separates all points in X , S has no common zero set in X , and if we are working with complex-valued functions, S is closed under conjugation.

Example. Let G be a locally compact group, with Haar measure μ , and consider the space of functions $L^1(G)$, where multiplication is convolution,

$$(f * g)(x) = \int f(y)g(y^{-1}x)d\mu(y)$$

One can then verify that

$$\|f * g\|_{L^1(G)} \leq \|f\|_{L^1(G)}\|g\|_{L^1(G)}.$$

Thus $L^1(G)$ is a Banach algebra, commutative when G is commutative. It does not always possess a unit, e.g. $L^1(\mathbf{R})$ is non-unital, but $L^1(\mathbf{Z})$ is unital. But we can always enlarge the space so it becomes a unital Banach algebra. Let $M(G)$ be the space of complex-valued finite Borel measures on G , and define convolution on $M(G)$ by letting

$$\int f d(\eta * \nu) = \int \int f(xy) d\eta(x) d\nu(y),$$

then $M(G)$ forms a unital Banach algebra, and the Dirac delta function δ_0 at the identity element of G is then a unit for $M(G)$. Thus one can consider the set

$$L^1(G)^\# = \{\delta_0 + f : f \in L^1(G)\},$$

which is the smallest unital Banach algebra containing $L^1(G)$ in $M(G)$.

The space $M(\mathbf{Z}) = L^1(\mathbf{Z})$ is a fairly concrete example. It can be taken as the space of integer valued sequences $c = \{c(n)\}$ (where $\sum |c(n)| < \infty$), with convolution

$$(a * b)(n) = \sum_{k \in \mathbf{Z}} a(k)b(n - k)$$

The Dirac delta function here is just the sequence δ_0 , where

$$\delta_i(k) = \begin{cases} 1 & k = i \\ 0 & \text{elsewise} \end{cases}$$

Since $\delta_i = (\delta^1)^i$, $L^1(\mathbf{Z})$ is generated by δ^1 and it's inverse.

One of the prime reasons to study Banach algebras is to study spaces of operators on a Banach space, and this is our primary example of a non-commutative Banach algebra.

Example. Let E be a Banach space. The space $B(E)$ of all bounded linear operators from E to itself is a Banach algebra with respect to the operator norm. It is a unital algebra, since it possesses the identity operator. The subset $K(E)$ of compact linear operators is a closed (double-sided) ideal of $B(E)$, and so is also a Banach algebra. $K(E)$ is unital if and only if E is finite dimensional.

Really, all that distinguishes a Banach space from a Banach algebra is a continuous multiplication structure, modulo the norm we use to define the topology of the space.

Proposition 7.1. *Let X be a Banach space equipped with a continuous multiplication structure. Then there is an equivalent norm on X which makes the space into a Banach algebra.*

Proof. Embed X in $B(X)$ by defining $T_M(N) = MN$ (since multiplication on the right is continuous, T_M truly is in $B(X)$ rather than just being a linear map). If $M_i \rightarrow M$, and $T_{M_i} \rightarrow T$, then we claim $T_M = T$. Indeed, since multiplication is continuous on the left,

$$T(N) = \lim_{i \rightarrow \infty} T_{M_i}(N) = \lim_{i \rightarrow \infty} M_i N = MN = T_M(N).$$

Thus the closed graph theorem implies that $M \rightarrow T_M$ is a continuous embedding of X in $B(X)$. If we define a norm $\|\cdot\|_A$ on X by setting $\|M\|_A = \|T_M\|$, then it follows that $\|M\|_A \sim \|M\|$, and $\|MN\|_A \leq \|M\|_A \|N\|_A$, so X has a Banach algebra structure with the norm $\|\cdot\|_A$. \square

The class of non-unital Banach algebras is not often that different from the class of unital Banach algebras. If A is a Banach algebra without unit, then we can always consider its unitization

$$A^\# = \{x + \lambda : x \in A, \lambda \in K\}$$

which will also be a Banach algebra by setting

$$(x_1 + \lambda_1)(x_2 + \lambda_2) = (x_1 x_2 + \lambda_1 x_2 + \lambda_2 x_1) + \lambda_1 \lambda_2$$

and defining

$$\|x + \lambda\| = \|x\| + |\lambda|.$$

A statement is often true about A if and only if it is true about $A^\#$, so that it suffices to solely study unital Banach algebras.

As we saw in the examples above, if $A = C_0(X)$, then we can identify $A^\#$ with $C(X \cup \{\infty\})$, where $X \cup \{\infty\}$ is the one-point compactification of X . There is a general school of Banach algebra theory known as *non commutative topology*, which studies Banach algebras A heuristically as if they were rings of continuous functions on some ‘non-commutative space’. Thus one can think of the unitization of a non-unital algebra as the ‘one-point compactification’ of this underlying non commutative space.

7.1 Involutions

An *involution* on a Banach algebra A is an antilinear map $M \mapsto M^*$, which satisfies

$$M^{**} = M \quad (MN)^* = N^*M^*.$$

One can think of an involution as the generalization of the conjugation operator $z \mapsto \bar{z}$ on \mathbf{C} , or of taking the adjoint of an operator on a Hilbert space. A *Banach $*$ Algebra* is a Banach algebra with a fixed involution. A *C^* algebra* is a Banach $*$ algebra A which satisfies

$$\|M^*M\| = \|M\|^2$$

for all $M \in A$.

An arbitrary involution on a Banach algebra can be rather discontinuous, which makes the general study of Banach $*$ algebras not too interesting. But most involutions in Banach $*$ algebra theory turn out to be continuous, in particular, often being isometries, i.e. satisfying $\|M^*\| = \|M\|$. This is true of any C^* algebra A , since if $M \in A$, then

$$\|M\|^2 = \|M^*M\| \leq \|M^*\| \|M\|$$

implies $\|M\| \leq \|M^*\|$, and

$$\|M^*\|^2 = \|MM^*\| \leq \|M^*\| \|M\|$$

implies $\|M^*\| \leq \|M\|$. The next result shows that any involution on a commutative, semisimple Banach algebra is continuous (an algebra A is semisimple if the intersection of all maximal ideals is trivial).

Proposition 7.2. *Let A be a commutative, semisimple Banach algebra. Then every involution on A is continuous.*

Proof. We rely on a fact (to be proved later), that every algebra homomorphism $\phi : A \rightarrow K$ is *automatically continuous*. The semisimplicity of A implies that for any $M_1, M_2 \in A$, $M_1 = M_2$ if and only if for every such algebra homomorphism $\phi : A \rightarrow K$, $\phi(M_1) = \phi(M_2)$. Now if $T : A \rightarrow A$ is an involution on A , and $\phi : A \rightarrow K$ is a homomorphism, then $\overline{\phi \circ T} : A \rightarrow K$ is also an algebra homomorphism, and thus continuous. Thus if $M_i \rightarrow M$, and $TM_i \rightarrow N$, then the continuity of $\overline{\phi \circ T}$ implies that

$$\overline{\phi(TM)} = \overline{\phi(N)},$$

i.e. $\phi(TM) = \phi(N)$. Thus $TM - N$ is in the kernel of every character ϕ . Since A is semisimple, this implies that $TM = N$. Thus by the closed graph theorem, T is continuous. \square

Remark. A similar proof shows that every homomorphism $T : A \rightarrow B$ from a Banach algebra A into a semisimple commutative Banach algebra B is continuous.

Example. Let H be a Hilbert space. Then for any bounded operator $T : H \rightarrow H$, we define the adjoint operator $T^* : H \rightarrow H$ by the equation

$$\langle Tx, y \rangle = \langle x, T^*y \rangle.$$

Now if $\|x\|, \|y\| \leq 1$, then

$$|\langle T^*Tx, y \rangle| = |\langle Tx, Ty \rangle| \leq \|T\|^2$$

and this is attained, for if there are $\|x_i\| \leq 1$ such that $\|Tx_i\| \rightarrow \|T\|$, then

$$\langle T^*Tx_i, x_i \rangle = \|Tx_i\|^2 \rightarrow \|T\|^2.$$

Thus $\|T^*T\| = \|T\|^2$, and so $B(H)$ is a C^* algebra. If we are working over \mathbf{C}^n , then each element of $B(\mathbf{C}^n)$ can be expressed as a matrix in $M_n(\mathbf{C})$, and the involution in $B(\mathbf{C}^n)$ is then equivalent to taking the complex transpose on sets of matrices, which gives a finite dimensional example. In particular, $M_n(\mathbf{C})$ is also a C^* algebra.

Example. The pointwise conjugation map $f \mapsto \bar{f}$ is an involution on $L^\infty(X)$, and for any $f \in L^\infty(X)$,

$$\|f\bar{f}\|_\infty = \| |f|^2 \|_\infty = \|f\|_\infty^2.$$

Thus $L^\infty(X)$ is a C^* algebra. Similarly, $C_0(X)$ and $C_b(X)$ are C^* algebras with the same involution. The central result of the study of commutative C^* algebras will show that all such algebras are isometric to $C_0(X)$ for some X .

Example. For $f \in A(\mathbf{D})$, if we define $f^*(z) = \overline{f(\bar{z})}$, then $*$ is an isometric involution on $A(\mathbf{D})$. But $A(\mathbf{D})$ is not a C^* algebra with this involution. To show this, let $f(z) = e^{iz}$. Then

$$(ff^*)(z) = e^{iz}\overline{e^{i\bar{z}}} = e^{iz}e^{-iz} = 1$$

Thus $\|ff^*\|_\infty = 1$, yet $f(-i) = f^*(i) = e > 1$, so

$$\|f\|_\infty \|f^*\|_\infty > \|ff^*\|_\infty.$$

Example. For any locally compact group G , the convolution algebra $L^1(G)$ has an involution $f^*(g) = \overline{f(g^{-1})}$. It also has the involution $\bar{f}(g) = \overline{f(g)}$. The first one turns out to be much more important to the theory of harmonic analysis on such groups. But neither give $L^1(G)$ the structure of a C^* algebra. As a characteristic example, if $G = \mathbf{Z}$, if $f = \delta_1 + \delta_0 - \delta_{-1}$, then

$$\bar{f} * f = f * f = \delta_2 + 2\delta_1 - \delta_0 - 2\delta_{-1} + \delta_{-2}$$

and

$$f * f^* = -\delta_2 + 3\delta_0 - \delta_{-2}.$$

Thus $\|f\|_{L^1(\mathbf{Z})} = 3$, whereas $\|\bar{f} * f\|_{L^1(\mathbf{Z})} = 7$ and $\|f * f^*\|_{L^1(\mathbf{Z})} = 5$, and neither of these quantities are equal to $\|f\|_{L^1(\mathbf{Z})}^2 = 9$. Similar examples can be constructed on general locally compact groups by using bump functions.

The next result is clearly not true in an arbitrary Banach algebra (e.g. by considering the trivial Banach algebra structure on any Banach space X).

Proposition 7.3. Let A be a C^* algebra. Then for any $M \in A$,

$$\|M\| = \sup\{\|MN\| : \|N\| \leq 1\} = \sup\{\|NM\| : \|N\| \leq 1\}$$

Proof. If $N = M^* \|M\|^{-1}$, $\|N\| = 1$, and

$$\|MN\| = \|MM^* \|M\|^{-1}\| = \|M\|$$

Thus the suprema on the right is greater than or equal to $\|M\|$. But if $\|N\| \leq 1$, then

$$\|MN\| \leq \|M\| \|N\| \leq \|M\|$$

so the suprema is also less than or equal to $\|M\|$. □

It follows that we may isometrically embed any C^* -algebra A into $B(A)$, by either considering the left or right multiplications. This is known as the *left or right regular representation*. Thus if A is a C^* algebra without identity, then we may embed A in $B(A)$, and then consider the smallest closed C^* -subalgebra of $B(A)$ which contains A and the identity. The general representation theory of C^* algebras, to be returned to at a later time. In particular, we will learn about the Gelfand Naimark Segal theorem, which

gives rise to a universal representation of A as a subalgebra of operators on a Hilbert space.

Given a non-unital Banach $*$ algebra A , there is a natural involution on $A^\#$ obtain by defining

$$(\lambda + M)^* = \overline{\lambda} + M^*$$

If A is a C^* algebra, so is $A^\#$, provided that we give it the norm

$$\|\lambda + M\| = \sup\{\|(M + \lambda)N\| : N \in A, \|N\| \leq 1\},$$

i.e. the norm of $\lambda + M$ is the operator norm of $\lambda + M$ obtained by embedding $A^\#$ in $B(A)$ by the left regular representation. The above proposition verifies that the embedding of A in $A^\#$ is an isometry under this norm.

The theory of C^* algebras is closely related to operator theory. The Gelfand Naimark Segal construction shows that all C^* algebras are isometric to some C^* subalgebra of bounded operators on a Hilbert space. Thus it makes sense to extend terminology from Hilbert space theory to the C^* domain. We say $M \in A$ is *self-adjoint* or *hermitian* if $M^* = M$. The set of all such elements forms a real subalgebra of A , analogous to the inclusion of the real numbers in the complex numbers. Here is a result which provides further evidence that this analogy is apt.

Proposition 7.4. *In any Banach $*$ algebra A , there is a unique way to write any $M \in A$ as $M = T + iS$, where T and S are self-adjoint elements of A .*

Proof. Write

$$T = \frac{M + M^*}{2} \quad S = \frac{M - M^*}{2i}$$

Then $T + iS = M$, and they are trivially verified to be self-adjoint. Now suppose $M = N + iL$, where N and L are self-adjoint. Then

$$0 = (N - T) + i(L - S)$$

Taking adjoints on both sides thus gives

$$(N - T) + i(L - S) = (N - T) - i(L - S)$$

And so $i(L - S) = -i(L - S)$, which can only occur when $L - S = 0$. But then $N - T = 0$ also. \square

An element M of a Banach $*$ algebra is *normal* if

$$M^*M = MM^*$$

The importance of this property results because if M is normal, then the smallest involutive subalgebra of A containing M will be commutative. An element M of a Banach $*$ algebra is *unitary* if

$$MM^* = M^*M = 1,$$

and the set of all such elements forms a closed subgroup $U(A)$ of $GL(A)$. A *projection* in a Banach $*$ algebra is an element P such that $P^2 = P$.

Chapter 8

Basic Spectral Theory of Banach Algebras

We shall begin our study into Banach algebras by analyzing criteria for invertibility, which coincides with the spectral theory of these algebras. For obvious reasons, we restrict our attention to unital algebras, though one can extend all concepts involved to a non-unital algebra A by embedding A in $A^\#$. The set of all units in a unital algebra A will be denoted $GL(A)$, and called the general linear group of A . The *spectrum* and *resolvent* of an element M of a Banach algebra A are defined respectively as

$$\sigma_A(M) = \{\lambda \in \mathbf{C} : \lambda - M \notin GL(A)\}$$

and

$$\rho_A(M) = \{\lambda \in \mathbf{C} : \lambda - M \in GL(A)\}$$

The resolvent is the complement of the spectrum. When the underlying algebra is canonical, we just use $\sigma(M)$ and $\rho(M)$ to denote the spectrum.

Example. Let X be a compact topological space, and consider $f \in C(X)$. Then $\sigma(f) = f(X)$, because the function

$$\frac{1}{\lambda - f(x)}$$

is well defined, bounded, and continuous if and only if $\lambda \notin f(X)$. If Ω is a domain in \mathbf{C} , then $\sigma_{A(\Omega)}(f) = \sigma_{C(\Omega)}(f)$, since under the assumptions of this example, the function $(\lambda - f(x))^{-1}$ will also be an analytic function. Similarly,

if X is a measure space, and $f \in L^\infty(X)$, then $\sigma(f) = \text{ess supp}(f)$, the essential support of f .

Example. Consider the algebra of $n \times n$ matrices M_n . An $n \times n$ matrix is invertible if and only if it is injective, so for an $n \times n$ matrix M , $\sigma(M)$ is precisely the set of eigenvalues of M . On the other hand, in a Banach space E , it is no longer true that a bounded operator $T : E \rightarrow E$ is invertible if and only if it is injective. Thus the point spectra $\sigma_p(T)$ of eigenvalues for T is a proper subset of the spectrum $\sigma(T)$.

The next lemma is incredibly important, and is an extension of the power series formula

$$\frac{1}{1-z} = \sum_{k=0}^{\infty} z^k$$

whose conclusion relies solely on the product and addition structure on the complex numbers, which all Banach algebras possess.

Lemma 8.1 (Neumann Series). *Let A be a Banach algebra. If $\|M\| < 1$, then $1 - M \in GL(A)$, and*

$$(1 - M)^{-1} = \sum_{k=0}^{\infty} M^k.$$

Proof. The right side converges absolutely by the comparison test, since $\|M^k\| \leq \|M\|^k$, and A is complete. This justifies the manipulation

$$(1 - M) \sum_{k=0}^{\infty} M^k = \sum_{k=0}^{\infty} (1 - M)M^k = \sum_{k=0}^{\infty} M^k - M^{k+1} = \lim_{n \rightarrow \infty} 1 - M^{n+1}$$

As $n \rightarrow \infty$, $M^{n+1} \rightarrow 0$, so the limit above tends to one. But $(1 - M)$ commutes with $\sum_{k=0}^{\infty} M^k$ (work with limits), since $\sum_{k=0}^{\infty} M^k$ is both a left and right inverse for $1 - M$. \square

Corollary 8.2. *Let A be a Banach algebra. If $M \in A$ and $\|1 - M\| < 1$, then $M \in GL(A)$, and*

$$M^{-1} = \sum_{k=0}^{\infty} (1 - M)^k$$

Corollary 8.3. *For any Banach algebra A , $GL(A)$ is an open subset of A .*

Proof. If $M \in GL(A)$, and if $\|M - N\| < 1/\|M^{-1}\|$, then

$$\|1 - M^{-1}N\| \leq \|M^{-1}\|\|M - N\| < 1$$

so $M^{-1}N \in GL(A)$, and thus $N \in GL(A)$. \square

Corollary 8.4. $\sigma(M)$ is a closed and bounded subset of \mathbf{C} , and $\rho(M)$ is open

Proof. The map $f : \lambda \mapsto \lambda - M$ is a continuous operation, for

$$\|(\lambda - M) - (\mu - M)\| = \|\lambda - \mu\| = |\lambda - \mu|$$

Since $GL(A)$ is open, $f^{-1}(GL(A)) = \rho(M)$ is open, hence $\sigma(M)$ is closed. If $|\lambda| > \|M\|$, then $\|M/\lambda\| < 1$, so $(1 - M/\lambda) \in GL(A)$, which means $\lambda - M$ is also invertible. Thus $\sigma(A)$ is closed and bounded, hence compact. \square

If we define the *spectral radius* $r(M)$ to be the maximum modulus of an element of $\sigma(M)$, then the last corollary proves that $r(M) \leq \|M\|$, though we need not have equality here (take a trivial Banach algebra X and equip it with a unit, in which it is then true that $\sigma(x) = \{0\}$ for each $x \in X$).

We shall see that the spectra of complex Banach algebras are never empty. This is why we mainly study complex algebras, rather than real algebras; there are even finite dimensional real operators with empty spectra.

Example. Consider the matrix

$$M = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}$$

as an element of the Banach algebra $M_2(\mathbf{R})$. The characteristic polynomial is calculated to be $\lambda^2 + 1$, so the matrix has no eigenvalues in the real numbers, and correspondingly, $\sigma(M) = \emptyset$. Over $M_2(\mathbf{C})$, we find $\sigma(M) = \{i, -i\}$, so the spectrum is non-empty in the complex extension of the Banach algebra.

Lemma 8.5. Inversion in an operator algebra A is continuous.

Proof. Let (M_n) be a sequence in $GL(A)$ converging to an invertible element M . Then, by continuity, $M_n M^{-1} \rightarrow 1$. If $\|1 - M_n M^{-1}\| < 1/2$, then

$$M_n^{-1} M = (M^{-1} M_n)^{-1} = \sum_{k=0}^{\infty} (1 - M^{-1} M_n)^k$$

It follows that

$$\begin{aligned}\|M_n^{-1}\| &\leq \|M^{-1}\| \|M_n^{-1}M\| \leq \|M^{-1}\| \sum_{k=0}^{\infty} \|(1 - M^{-1}M_n)^k\| \\ &\leq \|M^{-1}\| \sum 2^{-k} = 2\|M^{-1}\|\end{aligned}$$

Finally, we obtain convergence of inverses,

$$\|M_n^{-1} - M^{-1}\| = \|M_n^{-1}(M - M_n)M^{-1}\| \leq \|M_n^{-1}\| \|M - M_n\| \|M^{-1}\|$$

which tends to zero, since the first and third values are bounded, and the second converges to zero. Thus inversion is continuous. \square

The invertible elements $GL(A)$ cannot form a closed set of A , for A is path-connected, so there is a sequence of invertible elements x_i which converge to a non-invertible element. Fortunately, Banach algebras have a way of telling us when things are about to go wrong; the norm of the operator blows up. If A was an algebra of operators on a Banach space, one could prove this result from the uniform boundedness principle, and despite the fact we could prove the general result by embedding any Banach algebra into a space of operators of this form by taking an equivalent norm, we work more abstractly.

Lemma 8.6. *Let A be a Banach algebra containing a convergent sequence $M_i \rightarrow M$, such that $M_i \in GL(A)$ for each i , and $M \notin GL(A)$. Then $\|M_i^{-1}\| \rightarrow \infty$.*

Proof. If $\|M_n^{-1}\| \leq C$ for all n , then

$$\|1 - M_n^{-1}M\| \leq \|M_n^{-1}\| \|M_n - M\| \leq C\|M_n - M\|.$$

If we choose n large enough that $\|M_n - M\| < 1/C$, then $\|1 - M_n^{-1}M\| < 1$, which means $M_n^{-1}M \in GL(A)$, and thus $M \in GL(A)$. \square

Theorem 8.7. *If B is a closed subalgebra of A , and $M \in GL(B)$, then the entire connected component of M in $B \cap GL(A)$ is contained in $GL(B)$.*

Proof. $GL(A)$ is open, so $B \cap GL(A)$ is open in B . So too is $GL(B)$. If $M_n \rightarrow M$, with $M_n \in GL(B)$, and $M \in GL(A)$, then $\|M_n^{-1}\| \rightarrow \|M^{-1}\| < \infty$, so by the last lemma, M must be invertible in B . Thus we have argued that $GL(B)$ is a closed subset of $B \cap GL(A)$. But it is also open, and therefore must be a union of connected components of $B \cap GL(A)$. \square

Corollary 8.8. *If B is a closed subalgebra of A , then $\sigma_B(M)$ is obtained from $\sigma_A(M)$ by adding certain components of $\rho_A(M)$.*

Corollary 8.9. *If $A \subset B$ are algebras, and $\sigma_B(M) \subset \mathbf{R}$, then $\sigma_A(M) = \sigma_B(M)$.*

Proof. $\sigma_B(M)$ is a bounded subset of \mathbf{R} , so $\rho_B(M)$ is connected in \mathbf{C} . Hence $\sigma_A(M) = \sigma_B(M)$, since $\sigma_A(M) \neq \mathbf{C}$, and if we added any more points to the spectrum we would have to add the entire component, which would have to be all of $\rho_B(M)$. \square

8.1 Resolvents

The fundamental theorem of spectral theory relies on heavy complex analysis, hence its restricted application to Banach algebras over the complex numbers. Define the resolvent of M , defined on $\rho(M)$ by

$$R(z; M) = (z - M)^{-1}$$

Fixing M , we obtain a function $R : \rho(M) \rightarrow A$. Formally speaking, we should expect R to be an analytic function with derivative $-(z - M)^{-2}$. This turns out to be the case.

Lemma 8.10. *$R : \rho(M) \rightarrow A$ is analytic, in the sense that*

$$R'(z) = \lim_{w \rightarrow 0} \frac{R(z + w) - R(z)}{w}$$

is a well defined, convergent limit for all $z \in \rho(M)$.

Proof. The proof is a pure computation.

$$\begin{aligned} & \lim_{w \rightarrow 0} \frac{R(z + w, M) - R(z, M)}{w} \\ &= \lim_{w \rightarrow 0} \frac{(z + w - M)^{-1} - (z - M)^{-1}}{w} \\ &= \lim_{w \rightarrow 0} (z + w - M)^{-1} \frac{(z - M) - (z + w - M)}{w} (z - M)^{-1} \\ &= \lim_{w \rightarrow 0} -(z + w - M)^{-1} (z - M)^{-1} \\ &= -(z - M)^{-2}. \end{aligned}$$

\square

In functional analysis, we call such a mapping *strongly analytic*. A *weakly analytic* function $f : \mathbf{C} \rightarrow X$ is a map such that for each $\lambda \in X^*$, $\lambda \circ f$ is analytic. These notions are equivalent if X is a Banach space (proved via methods such as the Banach-Steinhaus theorem). Applying Louiville's theorem to the spectrum, we shall now prove that the spectrum of any element of a Banach algebra is non-empty.

Lemma 8.11. *A bounded weakly analytic function $f : \mathbf{C} \rightarrow X$ is constant.*

Proof. Let $\phi \in X^*$. Then $\phi \circ f$ is an analytic function, and

$$|(\phi \circ f)(z)| \leq \|\phi\| \|f\|_\infty$$

Louiville tells us there is $w \in \mathbf{C}$ such that $\phi \circ f = w$. Fix $x, y \in \mathbf{C}$. Then $\phi[f(x) - f(y)] = 0$ for all $\phi \in X^*$. By the Hahn Banach theorem, since ϕ was arbitrary, we must have $f(x) - f(y) = 0$, so $f(x) = f(y)$. Thus f is a constant function. \square

There is a deep relationship between complex analysis and the spectral theory of Banach algebras. We shall return to 'holomorphic functional analysis' later.

Theorem 8.12 (Gelfand-Mazur). *If A is a complex Banach algebra, then for any $M \in A$, $\sigma(M) \neq \emptyset$.*

Proof. Assume $\sigma(M)$ is empty. Then $R(z)$ is an analytic, entire function of z , and tends to zero at infinity, since if $|z| > N\|M\|$,

$$\begin{aligned} \|R(z)\| &= \|(z - M)^{-1}\| = \frac{1}{|z|} \left\| \sum_{k=0}^{\infty} \frac{M^k}{z^k} \right\| \leq \sum_{k=0}^{\infty} \frac{\|M\|^k}{|z|^{k+1}} \\ &\leq \frac{1}{|z|} \sum_{k=0}^{\infty} \frac{1}{N^{k+1}} = \frac{1}{N(N-1)\|M\|} \end{aligned}$$

But this implies $R = 0$. But then

$$R(z) = (z - M)^{-1} = 0,$$

and this is clearly impossible for any particular value of z . \square

Remark. We note how non-constructive this theorem is. In particular, this reflects the fact that, for a given explicit element M of a Banach algebra, it is often very difficult to compute $\sigma(M)$ explicitly.

A cute little theorem arises from this property of Banach algebras that will surprisingly find great use in the analysis of commutative Banach algebras.

Corollary 8.13. *Every complex Banach division algebra is isometric to \mathbf{C} .*

Proof. In any unital Banach algebra A , $\mathbf{C} \cdot 1$ is isometric to \mathbf{C} , for

$$\|\lambda \cdot 1\| = |\lambda| \|1\| = |\lambda|$$

Let A be a complex division algebra, and fix $M \in A$. Pick some $\lambda \in \sigma(A)$. Then $\lambda - M \notin GL(A)$, hence $\lambda - M = 0$, i.e. $M = \lambda$. Thus $A = \mathbf{C} \cdot 1 \cong \mathbf{C}$. \square

The real case is much more complicated. There are three Banach division algebras over the real numbers, namely \mathbf{R} , \mathbf{C} , and \mathbf{Q} (the quaternions), and it is much more difficult to show that these are the only three.

Theorem 8.14. *If A is a Banach algebra for which M exists such that*

$$\|x\| \|y\| \leq M \|xy\|$$

for all $x, y \in \mathbf{C}$, then A is isometric to \mathbf{C} .

Proof. Let $y \in \partial GL(A)$. Then there are $y_i \rightarrow y$ with $y_i \in GL(A)$. Thus $\|y_i^{-1}\| \rightarrow \infty$. But since

$$\|y_i\| \|y_i^{-1}\| \leq M$$

we must have $\|y_i\| \rightarrow 0$, so $y = 0$. Suppose $z \notin GL(A)$. Consider the line between z and 1 . Consider

$$\lambda = \inf\{\gamma \in [0, 1] : \gamma + (1 - \gamma)z \in GL(A)\}$$

then $\gamma + (1 - \gamma)z \in \partial GL(A)$, so $\gamma + (1 - \gamma)z = 0$, and since $\gamma \neq 1$ (since $GL(A)$ is open), we conclude $z = \gamma/(1 - \gamma)$, and since z is not invertible, we must have $z = 0$. Thus A is a division algebra, so A is isometric to \mathbf{C} . \square

Theorem 8.15. *σ is ‘continuous’, in the sense that for any open set $U \subset \mathbf{C}$ containing $\sigma(M)$, there is an open neighbourhood V of M such that if $N \in V$, $\sigma(N) \subset U$.*

Proof. The map $\lambda \mapsto \|(\lambda - M)^{-1}\|$ is a bounded continuous function in U^c , with some bound K . If $\|N\| < 1/K$, and $\lambda \in U^c$, then

$$\lambda - (M + N) = (\lambda - M)(1 - (\lambda - M)^{-1}N)$$

is invertible, since $(\lambda - M)$ is invertible, since $\|(\lambda - M)^{-1}N\| < 1$. It follows that $\sigma(M + N) \subset U$. \square

8.2 Spectral Radii

The *spectral radius* of an element $M \in A$ is defined to be

$$r(M) = \sup\{|\lambda| : \lambda \in \sigma(M)\}$$

It bounds where to look for the spectrum of M occurs. What is amazing is that we can define the spectral radius without any algebraic reference; this is crazy, since if we enlarge our Banach algebra, more elements become invertible, and thus the spectrum shrinks in size. The radius formula says that there still exists points on the boundary of the spectrum.

Lemma 8.16. *Let $M \in A$, $n \in \mathbf{N}$. Then, if $\lambda \in \sigma(M)$, $\lambda^n \in \sigma(M^n)$.*

Proof. Suppose $\lambda \in \sigma(M)$. Then

$$\lambda^n - M^n = (\lambda - M) \left(\sum \lambda^{n-1-k} M^k \right) = \left(\sum \lambda^{n-1-k} M^k \right) (\lambda - M)$$

If $\lambda - M$ was invertible, then $\lambda^n - M^n$ would also be invertible. \square

Theorem 8.17 (Spectral Radius Theorem).

$$r(M) = \lim_{n \rightarrow \infty} \|M^n\|^{1/n}$$

Proof. If $\lambda \in \sigma(M)$, then $\lambda^n \in \sigma(M^n)$, and therefore by Neumann's lemma, $|\lambda|^n \leq \|M^n\|$. We conclude

$$|\lambda| \leq \|M^n\|^{1/n}$$

taking extrema,

$$r(M) = \sup_{\lambda \in \sigma(M)} |\lambda| \leq \liminf_{n \rightarrow \infty} \|M^n\|^{1/n}$$

Set $R = r(M)^{-1}$ (which can be ∞ , if $r(M) = 0$), and $r = \|M\|^{-1}$. Let λ be a complex number with modulus less than R . Then $1/|\lambda| > r(M)$, so $1 - \lambda M \in GL(A)$. If $\phi \in A^*$, define

$$f : \lambda \mapsto \langle \phi, (1 - \lambda M)^{-1} \rangle$$

Then f is holomorphic in the disk of radius R . If $|\lambda| < r$, then $\|\lambda M\| < 1$, $1 - \lambda M \in GL(A)$, and

$$\phi((1 - \lambda M)^{-1}) = \sum_{k=0}^{\infty} \phi(\lambda^k M^k)$$

power series expansions are unique, hence this expansion should work in the whole disk of radius R . But ϕ was arbitrary, so the sequence $\lambda^k M^k$ must be bounded, by Banach Steinhaus. If λ is fixed, then there is C such that

$$|\lambda^n| \|M^n\| \leq C$$

for all n , so

$$\|M^n\|^{1/n} \leq \frac{C^{1/n}}{|\lambda|}$$

Hence

$$\limsup_{n \rightarrow \infty} \|M^n\|^{1/n} \leq \lim_{n \rightarrow \infty} \frac{C^{1/n}}{\lambda} = \frac{1}{\lambda}$$

Letting $\lambda \rightarrow R$, we obtain that $\limsup \|M^k\|^{1/k} \leq r(a)$. We have shown

$$\liminf \|M^k\|^{1/k} \geq r(M) \geq \limsup \|M^k\|^{1/k}$$

from which the theorem follows. \square

Corollary 8.18. *The spectral radius of a Banach algebra element is invariant of which Banach algebra the element is in. If B is a Banach subalgebra of A , with $M \in B$, then $r_A(M) = r_B(a)$.*

Example. *We can isometrically embed $A(\mathbf{D})$ in $C(\mathbf{T})$ via the map $f \mapsto f|_{\mathbf{T}}$, so that we may view $A(\mathbf{D})$ as a closed subspace of $C(\mathbf{T})$. These properties follows simply from the maximum modulus principle. Let $z : \mathbf{T} \rightarrow \mathbf{C}$ be the identity map. Then*

$$\sigma_{A(\mathbf{D})}(z) = \mathbf{D} \supsetneq \mathbf{T} = \sigma_{C_0(\mathbf{T})}(z|_{\mathbf{T}})$$

while the spectrum are different, the spectral radius is the same. Note that enlarging the spectrum in this way is the only possible change one can have in the spectrum, since we have to add a connected component of the resolvent when enlarging the algebra, and the interior of the unit disk is the only bounded connected component of the resolvent.

Corollary 8.19. *If A is a C^* algebra, and $M \in A$ is normal, then $r(M) = \|M\|$.*

Proof. We calculate that

$$\|M^2\| = \|(M^*)^2 M^2\|^{1/2} = \|(M^* M)(M^* M)^*\|^{1/2} = \|M^* M\| = \|M\|^2.$$

More general, by induction we can prove that

$$\|M^{2^n}\| = \|M\|^{2^n}.$$

But this means that

$$r(M) = \lim_{n \rightarrow \infty} \|M^{2^n}\|^{1/2^n} = \|M\|.$$

Since $r(M)$ is defined completely algebraically, without any reference to the norm of A , this shows $\|M\|$ depends solely on the algebraic structure of A . \square

Corollary 8.20. *Every C^* algebra A has a unique norm making it a C^* algebra.*

Proof. For any $M \in A$, $M^* M$ is self-adjoint, so

$$\|M\|^2 = \|M^* M\| = r(M^* M).$$

Thus the norm of M is defined by the algebraic structure of A . \square

8.3 Bounded Approximate Identities

We may not have an identity in a non-unital Banach algebra, but we may be able to ‘approximate’ an identity in some sense. An *approximate identity* for an algebra A is a net $\{E_\alpha\}$ such that for any $N \in A$,

$$\lim E_\alpha N = N \quad \lim N E_\alpha = N$$

Nets for which one of these equations hold are known as *left* or *right* approximate identities. A *bounded approximate identity* is an approximate identity for which $\sup \|E_\alpha\| < \infty$.

Example. *Let X be locally compact and Hausdorff. The set of all compact subsets is a directed, exhausting set. Using Urysohn’s lemma, find f_K , for each compact K , such that $f_K|_K = 1$, and $\|f_K\|_\infty \leq 1$. Fix $g \in C_0(X)$. Pick K such that $|g| \leq \varepsilon$ outside of K . Then, in K , $|f_K g - g| = 0$, and outside of K , $|f_K g - g| < 2\varepsilon$. Thus $\|f_K g - g\|_\infty < 2\varepsilon$. Thus $\lim_{K \rightarrow \infty} f_K g = g$, and so $\{f_K\}$ is a bounded approximate identity in $C_0(X)$.*

Example. Consider the net $P_r(z) = \sum_{\mathbf{z}} r^{|\mathbf{n}|} z^n$ in the group algebra $L^1(S^1)$, for $0 < r < 1$. Given a continuous function g , it is well known that $(P_r * g) \rightarrow g$ uniformly. Since, for any $g \in L^1(S^1)$, $\|P_r * g\|_{L^1} \lesssim \|g\|_{L^1}$, uniformly in r , it follows that $g * P_r \rightarrow g$ in $L^1(S^1)$ for any $g \in L^1(S^1)$. Thus $\{P_r\}$ is a bounded approximate identity.

Example. Bounded approximate identities appear in advanced Banach space theory. For instance, they arise in the study of compact operators. Let X be a Banach space, and let $A(X)$ be the closure of the family of all bounded finite rank operators on X . Then $A(X)$ is a right ideal in $B(X)$ contained in $K(X)$, the subalgebra of compact operators on X . We say X has the approximation property if $K(X) = A(X)$, i.e. all compact operators can be approximated in the operator norm by finite rank operators. We say X has the bounded approximation property if there exists a family of uniformly bounded finite rank operators $\{E_\alpha\}$ converging pointwise to the identity map on X . Under these assumptions, it then follows that the family $\{E_\alpha\}$ actually converges uniformly on compact subsets of X , from which it follows that for any compact operator $T \in K(X)$, $E_\alpha T$ is a family of finite rank operators converging in the operator topology to T . Thus if X has the bounded approximation property, then X has the approximation property. If a Banach space is reflexive, it is a result of Grothendieck that a space has the bounded approximation property if and only if $A(X) = K(X)$, i.e. all compact operators on X can be approximated by finite rank operators. We claim that X has the bounded approximation property if and only if $A(X)$ has a bounded left approximate identity.

If X has the bounded approximation property, and $\{E_\alpha\}$ are the uniformly bounded, finite rank operators converging uniformly to the identity on compact subsets of X , then it is simple to see that for any finite rank operator $T : X \rightarrow X$, $E_\alpha T$ converges in the operator norm to T . Indeed, if Y is the finite dimensional image of T , and B_Y is the closed unit ball in Y , then B_Y is compact, so E_α converges to the identity uniformly on B_Y , so $E_\alpha|_Y$ converges in the norm topology to $\text{id}|_Y$, and thus means that

$$\|E_\alpha T - T\|_{B(X)} = \|(E_\alpha - 1)T\|_{B(X)} \leq \|E_\alpha - 1\|_{B(Y)} \|T\|_{B(X)} \rightarrow 0.$$

A density argument shows the same is true for any $T \in A(X)$, so $\{E_\alpha\}$ is a bounded approximate identity for $A(X)$.

Conversely, suppose $\{E_\alpha\}$ is a bounded left approximate identity for $A(X)$. Without loss of generality, by a density argument we may assume each E_α is of

finite rank. For each $x_0 \in X$, and $\phi_0 \in X^*$, the operator $x_0 \otimes \phi_0 : X \rightarrow X$ defined by setting

$$(x_0 \otimes \phi_0)(x) = \phi_0(x) \cdot x_0$$

is clearly a finite rank operator. Thus we conclude that $E_\alpha(x_0 \otimes \phi_0)$ converges to $x_0 \otimes \phi_0$ in $B(X)$. But this means that

$$\|\phi_0(x)\| \|E_\alpha x_0 - x_0\| = o(\|x\|).$$

But this clearly means that $\|E_\alpha x_0 - x_0\| \rightarrow 0$. Thus $\{E_\alpha\}$ is a uniformly bounded family of operators converging pointwise to the identity, so X has the bounded approximation property.

Finding a Banach space without the approximation property proved to be a very challenging problem in 20th century mathematics. For instance, if H is a Hilbert space, then $B(H)$ always has the bounded approximation property (consider the bounded approximate identity given by orthogonal projections onto finite dimensional subspaces), as do the classical sequence spaces c_0 and l_p . The problem of proving or disproving that all Banach spaces have the approximation property was proposed by Stanislaw Mazur in 1936, promising a reward of a live goose as a prize for the solution to this problem. Forty years later, the problem was solved by the Swedish mathematician Per Enflo in 1972, who finally recieved a live goose as a reward.

Example. Let X be a Banach space with the bounded approximation property. Let $\{T_\alpha\}$ be a bounded net. If S is a compact operator, then $S(\overline{B_X})$ is precompact, and

$$\|T_\alpha S - S\| = \sup\{\|T_\alpha Sx - Sx\| : x \in B_X\} \leq \sup\{\|T_\alpha y - y\| : y \in \overline{S(B_X)}\}$$

The right side is the supremum over a compact set, and since $\{T_\alpha\}$ uniformly tends to the identity on compact sets, $\|T_\alpha S - S\| \rightarrow 0$. This shows $\{T_\alpha\}$ is a bounded left approximate identity, and that $A(X) = K(X)$.

Now suppose $A(X)$ has a bounded left approximate identity $\{T_\alpha\}$. Without loss of generality, we may assume each T_α is of finite rank, because if we choose, for each α , a bounded net $W_{\alpha,\beta}$ of finite rank operators such that $\lim_\beta W_{\alpha,\beta} = T_\alpha$, then $\{W_{\alpha,\beta}\}$ is a bounded left approximate identity. For $x \in X$, $\phi \in X^*$, let $x \otimes \phi : X \rightarrow X$ be defined by $(x \otimes \phi)(y) = \phi(y)x$. Then $x \otimes \phi$ obviously has finite rank. Pick y for which $\langle \phi, y \rangle = 1$. Then

$$\|T_\alpha x - x\| = \|T_\alpha(x \otimes \phi)(y) - (x \otimes \phi)(y)\| \leq \|y\| \|T_\alpha(x \otimes \phi) - (x \otimes \phi)\| \rightarrow 0$$

So the nets converge pointwise, which implies convergence on compact sets. To summarize, X has the bounded approximation property if and only if $A(X)$ has a bounded left approximate identity if and only if $K(X)$ has a bounded left approximate identity belonging to $A(X)$.

Lemma 8.21. *If a space has the bounded left and right approximation properties then it has the two sided approximation property.*

Proof. Let $\{M_\alpha\}$ be a left approximation identity, and $\{N_\beta\}$ a right approximation identity. We contend $\{M_\alpha + N_\beta - M_\alpha N_\beta\}$ is a two sided approximator. The limits below certainly converge, and the iterated limits must therefore equal the convergent factor, which is

$$\begin{aligned}\lim_{\alpha} \lim_{\beta} L(M_\alpha + N_\beta - M_\alpha N_\beta) &= \lim_{\alpha} LM_\alpha + L - LM_\alpha = L \\ \lim_{\beta} \lim_{\alpha} (M_\alpha + N_\beta - M_\alpha N_\beta)L &= \lim_{\beta} L + N_\beta L - N_\beta L = L\end{aligned}$$

So we have a two sided approximator. \square

Before the logician Paul Cohen got into logic, he was a functional analyst who contributed to the theory of approximation identities. We shall prove the theorem he contributed to the field, generalized to work over arbitrary modules. Let A be a Banach algebra. A *left Banach A -module* is a Banach space X , which is a module over A , such that for any $M \in A$, and $x \in X$,

$$\|Mx\| \leq \|M\| \|x\|$$

This makes the module operations continuous. As with Banach algebras, if module operations are continuous, then a Banach space can be made into a Banach A -module.

Now if A is a unital Banach algebra, and X is a left Banach A -module, then every $x \in X$ can be trivially written as My for some $M \in A$ (just choose M to be the identity). If A is a non-unital Banach algebra, and $x \in \overline{AX}$, then we can write x approximately as My . Cohen found that in an algebra with a bounded left approximate identity, we have $AX = \overline{AX}$, i.e. we can write x precisely as My for some $M \in A$ and $y \in X$.

Theorem 8.22 (Cohen's Factorization Theorem). *Let A be a Banach algebra with a bounded left approximation identity $\{E_\alpha\}$, bounded by some quantity $R > 0$. If X is a Banach A -module, let $x \in \overline{AX}$, and let $\varepsilon > 0$. Then there is $M \in A$ with $\|M\| \leq R$, $y \in \overline{AX}$ with $\|y - x\| < \varepsilon$, for which $x = My$.*

Proof. If A has an identity, then the proof is trivial. Since \overline{AX} is a closed subset of X , we might as well assume $\overline{AX} = X$. We may extend X to be an module over $A^\#$, by defining

$$(M + \lambda)x = Mx + \lambda x$$

Pick $\lambda > 0$ small enough such that

$$0 < \frac{\lambda}{1 - \lambda}K < 1$$

Define a net $\{L_\alpha\}$ in $A^\#$ by letting

$$L_\alpha = \lambda E_\alpha + (1 - \lambda) = (1 - \lambda) \left(1 + \frac{\lambda}{1 - \lambda} E_\alpha \right)$$

Then $L_\alpha N \rightarrow N$ for all $N \in A$, so $\{L_\alpha\}$ is a BLAI, and each L_α is invertible in $A^\#$ by the choice of λ , so

$$L_\alpha^{-1} = \frac{1}{1 - \lambda} \sum_{k=0}^{\infty} \left(-\frac{\lambda}{1 - \lambda} L_\alpha \right)^k$$

which implies

$$\begin{aligned} \|L_\alpha^{-1}N - N\| &= \|L_\alpha^{-1}N - L_\alpha^{-1}L_\alpha N\| \leq \|L_\alpha^{-1}\| \|N - L_\alpha N\| \\ &\leq \left(\frac{1}{1 - \lambda} \sum_{k=0}^{\infty} \left(\frac{\lambda K}{1 - \lambda} \right)^k \right) \|N - L_\alpha N\| \rightarrow 0 \end{aligned}$$

Therefore L_α^{-1} is also a BLAI. Fix $x \in X$, and let

$$\delta < 1, \frac{\varepsilon}{2 + \|x\|}$$

We will inductively construct a sequence α_n of indices such that the sequence $\{L_{\alpha_n} \dots L_{\alpha_1}\}$ and $\{L_{\alpha_1}^{-1} \dots L_{\alpha_n}^{-1}\}$ are convergent subnets. Choose α_1 such that

$$\|L_{\alpha_1}^{-1}x - x\| < \delta/2$$

If $\alpha_1, \dots, \alpha_n$ has been chosen. There is a unique element $T_n \in A$ such that

$$L_{\alpha_n} \dots L_{\alpha_1} = (1 - \lambda)^n + T_n$$

Pick α_{n+1} such that

$$\|L_{\alpha_{n+1}}x - x\| < \frac{\delta}{\|L_{\alpha_1}^{-1} \dots L_{\alpha_n}^{-1}\| 2^{n+2}} \quad \|L_{\alpha_{n+1}}T_n - T_n\| < 1/2^n$$

Then

$$\|L_{\alpha_1}^{-1} \dots L_{\alpha_n}^{-1} L_{\alpha_{n+1}}^{-1} x - L_{\alpha_1}^{-1} \dots L_{\alpha_n}^{-1} x\| \leq \|L_{\alpha_1}^{-1} \dots L_{\alpha_n}^{-1}\| \|x - L_{\alpha_{n+1}}^{-1} x\| < \frac{\delta}{2^{n+2}}$$

This implies $\lim_{n \rightarrow \infty} L_{\alpha_1}^{-1} \dots L_{\alpha_n}^{-1} x = y$ exists, and

$$\begin{aligned} \|y - x\| &= \left\| \sum_{n=0}^{\infty} (L_{\alpha_1}^{-1} \dots L_{\alpha_{n+1}}^{-1} x - L_{\alpha_1}^{-1} \dots L_{\alpha_n}^{-1} x) \right\| \\ &\leq \sum_{n=0}^{\infty} \|L_{\alpha_1}^{-1} \dots L_{\alpha_{n+1}}^{-1} x - L_{\alpha_1}^{-1} \dots L_{\alpha_n}^{-1} x\| < \sum_{n=0}^{\infty} \frac{\delta}{2^{n+2}} < \delta/2 \end{aligned}$$

By construction,

$$L_{\alpha_{n+1}}(L_{\alpha_n} \dots L_{\alpha_1}) = (1 - \lambda)^n L_{\alpha_{n+1}} + L_{\alpha_{n+1}} T_n$$

Thus

$$\begin{aligned} \|L_{\alpha_{n+1}} L_{\alpha_n} \dots L_{\alpha_1} - L_{\alpha_n} \dots L_{\alpha_1}\| &= \|(1 - \lambda)^n L_{\alpha_{n+1}} + L_{\alpha_{n+1}} T_n - T_n - (1 - \lambda)^n\| \\ &\leq (1 - \lambda)^n \|L_{\alpha_{n+1}} - 1\| + \frac{1}{2^n} \end{aligned}$$

Thus $L_{\alpha_n} \dots L_{\alpha_1}$ converge to some element M . Since A is closed in $A^\#$, $M \in A$, and moreover,

$$x = \lim_{n \rightarrow \infty} (L_{\alpha_1} \dots L_{\alpha_n})(L_{\alpha_1}^{-1} \dots L_{\alpha_n}^{-1} x) = My$$

Also, observe that

$$\begin{aligned} T_{n+1} &= L_{\alpha_{n+1}} \dots L_{\alpha_1} - (1 - \lambda)^{n+1} \\ &= L_{\alpha_{n+1}} T_n + L_{\alpha_{n+1}} (1 - \lambda)^n - (1 - \lambda)^{n+1} \\ &= L_{\alpha_{n+1}} T_n + (1 - \lambda)^n \lambda e_{\alpha_{n+1}} \end{aligned}$$

Thus

$$\|T_{n+1} - T_n\| = \|(L_{\alpha_{n+1}} T_n - T_n + (1 - \lambda)^n \lambda e_{\alpha_{n+1}})\| \leq 1/2^n + (1 - \lambda)^n \lambda K$$

Since $M = \lim T_n$,

$$\|M\| \leq \|T_1\| + \left(\sum_{n=1}^{\infty} (1-\lambda)^n \lambda K \right) + \delta \leq \lambda K + (1-\lambda)K + \delta = K + \delta$$

Let $M' = \frac{K}{K+\delta}M$, $y' = \frac{K+\delta}{K}y$. Then $M'y' = My = x$, and by the choice of δ (kept hidden all this time),

$$\|x - y'\| \leq \frac{\|Kx - Ky\| + \|\delta y\|}{K} < \delta + \delta\|x\| + \delta^2 \leq \delta(2 + \|x\|) = \varepsilon$$

And we have verified what was needed. \square

Corollary 8.23. *If A is a Banach algebra with a bounded left approximate identity, then each $M \in A$ may be written $M = NL$, with $N, L \in A$.*

Thus, for instance, for any locally compact group G , since $L^1(G)$ has an approximate identity, any $f \in L^1(G)$ can be written as $f = g * h$ for some $g, h \in L^1(G)$. Moreover, for any $\varepsilon > 0$, we can pick g and h such that $\|g\|_{L^1(G)} \leq 1$, and $\|f - h\|_{L^1(G)} < \varepsilon$, or vice versa.

Corollary 8.24. *Let A be a Banach algebra with a BLAI, and let $\{M_n\}$ be a sequence in A converging to 0. Then there is $M \in A$ and $\{N_m\}$ in A tending to zero such that $M_k = MN_k$.*

Proof. Let X be the set of all sequences in A that converge to zero, with the $\|\cdot\|_{\infty}$ norm. Then X is a left A module, and $X = \overline{AX}$, since A has a BLAI. Applying Cohen's theorem, we find that we may write

$$(M_1, M_2, \dots) = M(N_1, N_2, \dots)$$

hence $M_i = MN_i$ for each i . \square

8.4 The Holomorphic Functional Calculus

Given a polynomial $P(x) = \sum a_i x^i$, a Banach algebra A , and $M \in A$, it is natural to overload the definition of P by setting $P(M) = \sum a_i M^i$. This does not really require any topological info about A . In this section, we use this topology to extend this so that we can consider $f(M)$, for a large

family of analytic functions f . As a basic example of this, in matrix theory it is useful to define the exponential of $M \in A$ by

$$e^M = \sum_{k=0}^{\infty} \frac{M^k}{k!}$$

We may define $\sin(M)$ and $\cos(M)$ similarly. More generally, for any power series $f = \sum_{k=0}^{\infty} c_k(X - \alpha)^k$ with a radius of convergence R , we may define

$$f(M) = \sum_{k=0}^{\infty} c_k(M - \alpha)^k$$

And this function is defined for any M satisfying $\|M - \alpha\| < R$. For analytic functions not given in a power series representation, it is non-trivial to define them for elements of a Banach algebra, but very useful.

To do this, we utilize the Cauchy integral formula

$$f(w) = \frac{1}{2\pi i} \int_{\gamma} \frac{f(z)}{z - w} dz,$$

where γ is a simple positively oriented curve containing w in its interior. We use this formula to overload the function f . Given a Banach algebra A containing an element M , a domain $\Omega \subset \mathbb{C}$ containing $\sigma(M)$, an analytic function $f : \Omega \rightarrow A$, and a positively oriented curve γ containing $\sigma(M)$ in its interior, we define

$$f(M) = \frac{1}{2\pi i} \int_{\gamma} f(z)(z - M)^{-1} dz.$$

The right hand side is a vector-valued integral, whose integrand is smooth on its compact domain. Thus we can define this integral using Riemann sums (i.e. using the *Bochner integral*). We can reduce the study of such quantities to standard complex contour integrals by the fact that for any linear functional $\Lambda : A \rightarrow \mathbb{C}$,

$$\Lambda \left(\frac{1}{2\pi i} \int_{\gamma} f(z)(z - M)^{-1} dz \right) = \frac{1}{2\pi i} \int_{\gamma} \Lambda(f(z)(z - M)^{-1}) dz.$$

In particular, it follows from such a reduction that $f(M)$ is actually independent of the contour γ used. We have thus laid the foundations for the holomorphic functional calculus.

Theorem 8.25. *The map $f \mapsto f(M)$ is the unique algebra homomorphism from the algebra $\mathcal{O}(\sigma(M))$ into A extending the evaluation of polynomials, and which is continuous under the topology of locally uniform convergence.*

Proof. If the integrands of a Bochner integral converge uniformly, then the integrals converge. This gives continuity of the map $f \mapsto f(M)$. The map $f \mapsto f(M)$ is obviously linear. Next, we show that for any polynomial f , and any contour integral γ around $\sigma(M)$,

$$f(M) = \frac{1}{2\pi i} \int_{\gamma} f(z)(z - M)^{-1} dz.$$

By linearity, we may assume $f(z) = z^k$. We may assume our contour is large enough that $\|(z - M)^{-1}\| < 1$ on γ . Then

$$\begin{aligned} \frac{1}{2\pi i} \int_{\gamma} z^n (z - M)^{-1} dz &= \frac{1}{2\pi i} \int_{\gamma} z^{n-1} \sum_{k=0}^{\infty} \frac{M^k}{z^k} dz \\ &= \frac{1}{2\pi i} \sum_{k=0}^{\infty} M^k \int_{\gamma} z^{n-1-k} dz = M^n, \end{aligned}$$

since $\int_{\gamma} z^{n-1-k} dz \neq 0$ only when $n - 1 - k = -1$ ($n = k$). Now we need to show the multiplicity of the evaluation map. This is done by brute calculation. Let $f, g \in \mathcal{O}(D)$ be given. Pick γ be as required, and consider

another $\tilde{\gamma}$ satisfying $\text{Int } \tilde{\gamma} \supset \overline{\text{Int } \gamma}$, $\mathbf{C} - D \subset \text{ext } \tilde{\gamma}$. Then

$$\begin{aligned}
f(M)g(M) &= \frac{-1}{4\pi^2} \left(\int_{\gamma} f(z)(z-M)^{-1} dz \right) \left(\int_{\tilde{\gamma}} g(w)(w-M)^{-1} dw \right) \\
&= \frac{-1}{4\pi^2} \int_{\gamma} \int_{\tilde{\gamma}} f(z)g(w)(z-M)^{-1}(w-M)^{-1} dw dz \\
&= \frac{-1}{4\pi^2} \int_{\gamma} \int_{\tilde{\gamma}} f(z)g(w) \left(\frac{1}{w-z} ((z-M)^{-1} - (w-M)^{-1}) \right) dw dz \\
&= \frac{-1}{4\pi^2} \int_{\gamma} f(z) \left(\int_{\tilde{\gamma}} \frac{g(w)}{w-z} dw \right) (z-M)^{-1} dz \\
&\quad + \frac{1}{4\pi^2} \int_{\tilde{\gamma}} g(w) \left(\int_{\gamma} \frac{f(z)}{w-z} dz \right) (w-M)^{-1} dw \\
&= \frac{1}{2\pi i} \int_{\gamma} f(z) \left(\frac{1}{2\pi i} \int_{\tilde{\gamma}} \frac{g(w)}{w-z} dw \right) (z-M)^{-1} dz \\
&= \frac{1}{2\pi i} \int_{\gamma} f(z)g(z)(z-M)^{-1} dz \\
&= (fg)(M)
\end{aligned}$$

and thus we have an algebra homomorphism.

Finally, we prove uniqueness. This would be trivial if $\Omega = \mathbf{C}$, since all functions are limits of polynomials. More generally, we apply Runge's theorem, which says that for any compact set K , if S is a set containing at least one point from each bounded connected component of $\mathbf{C} - K$, then every analytic function on K is the uniform limit of a family of rational functions with poles solely in S . Now suppose $\Lambda : \mathcal{O}(\sigma(M)) \rightarrow A$ is another homomorphism satisfying the properties above. Then for any two polynomials P and Q , such that the zeroes of Q are disjoint from $\sigma(M)$, $1/Q \in \mathcal{O}(\sigma(M))$, and so $\Lambda(1/Q) = \Lambda(Q)^{-1}$. Thus we calculate that

$$\Lambda(P/Q) = \Lambda(P)\Lambda(1/Q) = \Lambda(P)\Lambda(Q)^{-1} = P(M)Q(M)^{-1} = (P/Q)(M).$$

Thus Λ agrees with the holomorphic functional calculus on rational functions, which shows $\Lambda(f) = f(M)$ for all $f \in \mathcal{O}(\sigma(M))$ by continuity and Runge's theorem. \square

Example. Let K be compact, and let $f \in C(K)$. Let Ω be an open neighbourhood of $\sigma(f) = f(K)$. Consider the map from $\mathcal{O}(\Omega) \rightarrow C(K)$ mapping g to $g \circ f$. This map satisfies the properties of the holomorphic functional calculus, so it really is the holomorphic calculus in disguise.

Theorem 8.26. Let $F : A \rightarrow B$ be a homomorphism between two Banach algebras, and suppose $M \in A$. Then $\sigma(F(M)) \subset \sigma(M)$, and if g is holomorphic in a neighbourhood of $\sigma(M)$, then $(g \circ F)(M) = (F \circ g)(M)$.

Proof. That $\sigma(f(M)) \subset \sigma(M)$ is trivial, for if $\alpha - M$ is invertible, then $f(\alpha - M) = \alpha - f(M)$ is invertible. If g is a polynomial, then the theorem follows from the multiplicative property. But then the theorem is true for limits of polynomials, and hence for all g . \square

The ultimatum of the holomorphic functional calculus is the proof of the spectral mapping theorem, that asserts that the ‘operation’ of the spectrum commutes with holomorphic functions. To extend this theorem to non-commutative algebras, we need an algebraic trick. We introduce the center $Z(A)$ of an algebra, which is defined exactly how it is defined in all other algebraic structures. It is a closed subalgebra of A , for if $M_i \rightarrow M$, and each $M_i \in A$, then

$$MN = \lim M_i N = \lim N M_i = N M$$

If S is a subset of A , then we may consider

$$Z(S) = \{M \in A : (\forall N \in S : MN = NM)\}$$

which is a commutative subalgebra of A . For any subset S , the double centralizer $Z(Z(S))$ also contains S . $Z(Z(S))$ is commutative if $S \subset Z(S)$, for then $Z(Z(S)) \subset Z(S)$, and thus commutes with itself by definition.

Now consider the commutative subalgebra $B = Z(Z(\{M\}))$ of an algebra A . If $\lambda - M$ is invertible in A , there is $N = (\lambda - M)^{-1}$. But then if $K \in Z(\{M\})$, then K commutes with $\lambda - M$, and

$$KN = N(\lambda - M)KN = NK(\lambda - M)N = NK$$

which implies that $\lambda - M$ is invertible in $Z(Z(\{M\}))$, so $\sigma_A(M) = \sigma_B(M)$.

Theorem 8.27 (Spectral Mapping Theorem). If $\sigma(M) \subset \Omega$, where Ω is open, and $f \in \mathcal{O}(\Omega)$, then $\sigma(f(M)) = f(\sigma(M))$.

Proof. First, suppose A is commutative. By Gelfand theory,

$$\sigma(f(M)) = \widehat{f(M)}(\Phi_A) = f(\widehat{M}(\Phi_A)) = f(\sigma(M))$$

If A is not commutative, consider the commutative subalgebra

$$B = Z(Z(\{M\})).$$

If f is analytic in a neighbourhood of $\sigma_A(M)$, then $f(M) \in B$, for there is surely a homomorphism from analytic functions in a neighbourhood of $\sigma_A(M)$ into B , and the embedding of B in A produces the unique homomorphism required. Thus we have justified the computation

$$\sigma_A(f(M)) = \sigma_B(f(M)) = f(\sigma_B(M)) = f(\sigma_A(M))$$

and the spectral mapping theorem is proved in general. \square

The holomorphic functional calculus is incredibly useful for it allows us to apply calculations on the complex plane to arbitrary algebras. Here are some immediate uses.

Theorem 8.28. *Suppose A is a Banach algebra, $M \in GL(A)$, and $\sigma(M)$ does not separate the origin from ∞ . Then*

- *For each m , there is N such that $M = N^m$. If $\sigma(M) \subset \mathbf{R}^+$, then we may choose the N such that $\sigma(N) \subset \mathbf{R}^+$.*
- *There is N such that $M = e^N$.*
- *if B is the closed subalgebra of A generated by M , then $M^{-1} \in B$.*

Proof. By assumption, we may choose a branch $f(z)$ of $z^{1/m}$ which is holomorphic on a neighborhood of $\sigma(M)$, and if $\sigma(M) \subset \mathbf{R}^+$, we may choose $f(z)$ to be the principal branch. Thus we can define $N = f(M)$. Since $f^m(z) = z$ for all z on a neighborhood of $\sigma(M)$, it follows that $N^m = f(M)^m = f^m(M) = M$. The spectral mapping theorem shows that $\sigma(N) = \sigma(f(M)) = f(\sigma(M))$, so if f is the principal branch, and $\sigma(M) \subset \mathbf{R}^+$, then $\sigma(N) \subset \mathbf{R}^+$. Similar calculations involving the logarithm allow us to define N such that $M = e^N$. Finally, by Runge's theorem, for each $\varepsilon > 0$, we can find a polynomial $f(z)$ such that $|f(z) - 1/z| \leq \varepsilon$ on $\sigma(M)$. Then $\|f(M) - M^{-1}\| \leq \varepsilon$. Taking $\varepsilon \rightarrow 0$ gives the final point. \square

One can apply this result, for instance, for unitary, or self adjoint elements, of C^* algebras.

Proposition 8.29. *In a C^* algebra, if M is unitary, then $\sigma(M) \subset \mathbf{T}$.*

Proof. Because M is normal, $r(M) = \|M\| = 1$. But $M^{-1} = M^*$ is also unitary, so $r(M^{-1}) = 1$, and by the spectral mapping theorem,

$$\sigma(M^{-1}) = \sigma(M)^{-1}$$

which implies that both spectra lie on \mathbf{T} . □

Proposition 8.30. *If M is self-adjoint, then $\sigma(M) \subset \mathbf{R}$.*

Proof. The operator e^{iM} is unitary, for

$$e^{-iM} = \sum_{k=0}^{\infty} \frac{(-i)^k M^k}{k!} = \sum_{k=0}^{\infty} \left(\frac{(iM)^k}{k!} \right)^* = (e^{iM})^*$$

The spectral mapping theorem implies that $\sigma(e^{iM}) = e^{i\sigma(M)} \subset \mathbf{T}$, which implies $\sigma(M) \subset \mathbf{R}$. □

Corollary 8.31. *Let A be a C^* algebra, and let B be a C^* subalgebra. Then $GL(B) = B \cap GL(A)$.*

Proof. Suppose $N \in B \cap GL(A)$. Then $N^* \in B \cap GL(A)$, for $(N^{-1})^* = (N^*)^{-1}$. Thus $NN^* \in B \cap GL(A)$. Moreover, NN^* is invertible, so $\sigma_A(NN^*) \subset \mathbf{R} - \{0\}$. Now $\sigma_B(NN^*)$ is obtained by adding bounded connected components of $\rho_A(NN^*)$ to $\sigma_A(NN^*)$, and since $\rho_A(NN^*)$ has no bounded connected components, it follows that $\sigma_B(NN^*) = \sigma_A(NN^*)$. But this means that $NN^* \in GL(B)$. But then N has a right inverse in B . By symmetry $N^*N \in GL(B)$, so N has a left inverse in B . But if N has a left and right inverse, then $N \in GL(B)$. □

Corollary 8.32. *If A is a C^* algebra, and B is a C^* -subalgebra, then for any $M \in B$, $\sigma_A(M) = \sigma_B(M)$.*

Corollary 8.33. *If A is a C^* algebra, and $M \in GL(A)$ is normal, then M^{-1} is approximable by expressions of the form*

$$N = \sum_{\alpha, \beta} a_{\alpha\beta} M^{\alpha} (M^*)^{\beta},$$

i.e. for any $\varepsilon > 0$, we can find N as above such that $\|M^{-1} - N\| \leq \varepsilon$.

Proof. Let B be the smallest C^* subalgebra of A containing M . Then the family of expressions of the form above is dense in B . By the result we have just proven, $M \in GL(B)$, and so the result immediately follows. \square

Even for finite dimensional matrix algebras, the holomorphic functional calculus has nontrivial repercussions.

Corollary 8.34. *If $M \in GL_n(\mathbf{C})$, then for each m , there is $N \in GL_n(\mathbf{C})$ such that $M = N^m$, and there is N such that $M = e^N$.*

Example. *Now consider the exponential function $\exp : A \rightarrow A$. If M and N commute, it is easy to see from the power series representation that*

$$\exp(M + N) = \exp(M)\exp(N).$$

In particular, this implies that $\exp(M)$ is invertible, and $\exp(-M) = \exp(M)^{-1}$. Thus for each M , the image of the map

$$t \mapsto \exp(tM)$$

is a one-parameter subgroup of $G(A)$, which gives us a path from the identity $\exp(0)$ to $\exp(M)$. The connected component of $G(A)$ containing the identity is the principal component, denoted $G_0(A)$. Since $N \mapsto MNM^{-1}$ is a continuous map leaving the identity fixed, it maps $G_0(A)$ to itself, so $G_0(A)$ is a normal subgroup. The quotient group $G(A)/G_0(A)$ is sometimes called the index group.

Theorem 8.35. *If $P(X) = (X - \alpha_1)^{m_1} \dots (X - \alpha_s)^{m_s}$, and $P(M) = 0$, then*

$$\sigma(M) \subset \{\alpha_1, \dots, \alpha_s\},$$

and if Ω is an open subset of the plane containing $\alpha_1, \dots, \alpha_s$, then for any $f \in \mathcal{O}(\Omega)$, there is a polynomial Q of degree less than $m_1 + \dots + m_s$, such that $f - Q = gP$, for some $g \in \mathcal{O}(\Omega)$, so $f(M) = Q(M)$.

Proof. Applying the spectral theorem, we find $\{0\} = \sigma(P(M)) = P(\sigma(M))$, from which the spectrum of M is restricted to lie in the roots of P . The Laurent series about each α_i gives constants c_{it} such that

$$g(z) = f(z)/P(z) - \sum_{i=1}^s \sum_{k=1}^{m_i} \frac{c_{ik}}{(z - \alpha_i)^k}$$

is holomorphic in Ω , and

$$gP = f - \sum_{i=1}^s \sum_{k=1}^{m_i} c_{ik} P(z - \alpha_i)^{-k}$$

constructing Q with degree less than n . \square

In non-unital algebras, we cannot evaluate any polynomial of the form $\sum a_i z^i$, since $a_0 \in \mathbf{C}$ does not necessarily have an interpretation in the algebra. If $a_0 = 0$, we can evaluate $\sum a_i z^i$, and thus any holomorphic function f , provided $f(0) = 0$. To verify this, one need only use the simple trick of switching from a non-unital algebra A to a unital algebra $A^\#$.

For a Banach space X , a weaker version of the spectral mapping theorem applies to the *point spectrum* $\sigma_p(T)$ of bounded operators $T : X \rightarrow X$.

Theorem 8.36. *Suppose $T : X \rightarrow X$ is a bounded operator, and $f \in \mathcal{O}(\sigma(T))$. Then $f(\sigma_p(T)) \subset \sigma_p(f(T))$, with equality if f is not constant on any connected component of the domain of f .*

Proof. Suppose $Tx = \lambda x$. The function $g(z) = (f(z) - f(\lambda))/(z - \lambda)$ is analytic, and so

$$[f(T)x - f(\lambda)x] = [g(T) \circ (T - \lambda)]x = g(T)0 = 0.$$

Thus $f(T)x = f(\lambda)x$. Now if $f(T)x = \gamma x$, and $f - \gamma$ does not vanish in any connected component of the domain of f , then $f(z) - \gamma$ has finitely many zeroes ξ_1, \dots, ξ_n , counted up to multiplicity. Thus the function

$$g(z) = \frac{(z - \xi_1) \dots (z - \xi_n)}{f(z) - \gamma}$$

is analytic and non-vanishing on the domain of f , and so

$$(T - \xi_1) \dots (T - \xi_n)x = g(T)[f(T) - \gamma]x = 0.$$

But this means that there is some nonzero y with $(T - \xi_i)y = 0$, and so $\xi_i \in \sigma_p(T)$. Since $f(\xi_i) = \gamma$, this shows that $f(\sigma_p(T)) = \sigma_p(f(T))$ if f is non constant on any connected component of its domain. \square

As mentioned briefly at the beginning of this section, given a non-unital algebra A , we extend the concepts of this chapter to elements of A by embedding A in $A^\#$. Thus, for instance, if $M \in A$, then

$$\sigma_A(M) = \{\lambda : \lambda - M \text{ is not invertible in } A^\#\}.$$

The spectral radius formula remains true in this setting, as does the result that $\sigma(M)$ is non-empty (though it is rather trivial in this setting, since $0 \in \sigma(M)$ for any $M \in A$). If $M \in A$, and f is holomorphic on a neighborhood of $\sigma(M)$, then we can define $f(M) \in A^\#$. Moreover, we have $f(M) \in A$ precisely when $f(0) = 0$, because if we consider the continuous linear functional $\Lambda(\lambda + M) = \lambda$ on $A^\#$, then $\Lambda(f(M)) = f(0)$ for any polynomial f , and we can apply a density argument. Thus we have a holomorphic functional calculus on A provided we consider only holomorphic functions vanishing at the origin.

Chapter 9

Gelfand Theory

Gelfand found a very powerful way of transforming commutative Banach algebras into Banach algebras of continuous functions over a compact space, which generalizes the spectral theory of commuting operators, as well as the Fourier transform. The ingenious trick to Gelfand theory is that one can recover the points in the set K via the algebraic structure of $C(K)$ by studying the maximal ideals of $C(K)$. Thus given a general commutative Banach algebra A , we can study the set K of maximal ideals of A , equipped this set with a topology, and then consider a natural homomorphism of A into $C(K)$ which has powerful algebraic properties. For reasons which might become clear from the last chapter (the main reason being that the only complex Banach division algebras are fields), in this chapter we solely work over algebras defined on the complex numbers.

9.1 Algebra Homomorphisms and Ideals

Given an algebra A containing an ideal \mathfrak{a} , we can consider the quotient ring A/\mathfrak{a} . We can consider left, right, and two-sided ideals, and we'll denote them by gaudy letters such as \mathfrak{a} and \mathfrak{b} . If \mathfrak{a} is a *closed* ideal in a Banach algebra, then A/\mathfrak{a} naturally has the structure of a Banach algebra, equipped with the norm defined, for any coset $S \in A/\mathfrak{a}$,

$$\|S\|_{A/\mathfrak{a}} = \inf_{s \in S} \|s\|_A.$$

Ideals are naturally connected with homomorphisms $f : A \rightarrow B$ between algebras. Every ideal is the kernel of some homomorphism, and the kernel

of every homomorphism is an ideal. Similarly, *closed ideals* in a Banach algebra correspond to continuous homomorphisms. The most tractable homomorphisms to study are the *characters*, homomorphisms from a ring A to \mathbb{C} . In this section, we will deduce their structure.

Lemma 9.1. *Let A be a unital Banach algebra and let $\phi : A \rightarrow \mathbb{C}$ be a character. Then $\|\phi\| \leq 1$.*

Proof. Let \mathfrak{a} be the kernel of ϕ . Then \mathfrak{a} cannot be all of A . Thus \mathfrak{a} does not contain any invertible elements of A . If $f(M) = \lambda$, then $\lambda - M \in \mathfrak{a}$, so $\lambda - M \in \sigma(M)$. This means $|\lambda| \leq r(M) \leq \|M\|$. \square

Corollary 9.2. *Every character is a continuous linear map.*

Lemma 9.3. *Every maximal ideal of a unital Banach algebra is closed.*

Proof. Let \mathfrak{a} be a maximal ideal of an algebra A . It is easy to show the closure of any ideal is an ideal. It follows that either $\bar{\mathfrak{a}} = \mathfrak{a}$ (so that \mathfrak{a} is closed), or \mathfrak{a} is dense in A . Suppose the second option holds. Let $M \in GL(A)$ be chosen. Then there is $M_i \in \mathfrak{a}$ converging to M . But then the M_i are eventually invertible, since $GL(A)$ is open, from which we conclude $\mathfrak{a} = A$, a contradiction. \square

Corollary 9.4. *If \mathfrak{a} is a maximal ideal in a commutative, unital Banach algebra A , then A/\mathfrak{a} is isometric to \mathbb{C} .*

Proof. For then A/\mathfrak{a} is a Banach division algebra. \square

Thus it follows that the kernel of every character on a commutative, unital Banach algebra A is a maximal ideal in A , and conversely, every maximal ideal corresponds to the kernel of some character. Moreover, this character is uniquely determined.

Lemma 9.5. *If $\phi_1 : A \rightarrow \mathbb{C}$ and $\phi_2 : A \rightarrow \mathbb{C}$ have the same kernel, then $\phi_1 = \phi_2$.*

Proof. Fix $M \in A$. Then

$$\phi_1(M - \phi_1(M)) = \phi_1(M) - \phi_1(M) = 0.$$

Thus $M - \phi_1(M)$ is in the kernel of ϕ_1 . But this means $M - \phi_1(M)$ is also in the kernel of ϕ_2 , so

$$\phi_2(M - \phi_1(M)) = \phi_2(M) - \phi_1(M) = 0.$$

Thus $\phi_1(M) = \phi_2(M)$. \square

We denote the set of characters of a Banach algebra by A by Φ_A , and call it the *character space* of A . Gelfand's theory is the study of Φ_A , and its relation to A , especially when A is a commutative Banach algebra.

Example. Let K be a compact space, and consider the algebra $C(K)$ of continuous functions. Given a closed subset C of K , the set

$$\{f \in C(K) : f|_C = 0\}$$

is an ideal of $C(K)$. We contend these are all such closed ideals. Let \mathfrak{a} be an arbitrary closed ideal of $C(K)$. It is easy to see that the set

$$C = \{x \in K : f(x) = 0 \text{ for all } f \in \mathfrak{a}\}.$$

are closed. Conversely, suppose $g \in C(K)$ and $g|_C = 0$. Fix $\varepsilon > 0$, and let $V = g^{-1}(B_\varepsilon(0))$. For each $y \in K - C$, there is a function $f_y \in \mathfrak{a}$ with $f_y(y) = g(y)$. Let $U_y = (f_y - g)^{-1}(B_\varepsilon(0))$. Then $\{V\} \cup \{U_y\}$ is an open cover of K , and thus has a finite subcover $V, U_{y_1}, \dots, U_{y_n}$. Let $h_V, h_{y_1}, \dots, h_{y_n}$ be a partition of unity subordinate to the cover. Consider the function $\tilde{g} = \sum h_{y_i} f_{y_i} \in \mathfrak{a}$. For $x \in K - V$,

$$|\tilde{g}(x) - g(x)| = \left| \sum h_{y_i}(x)(f_{y_i}(x) - g(x)) \right| \leq \sum h_{y_i}(x)|f_{y_i}(x) - g(x)| < \varepsilon$$

For $x \in V$,

$$|\tilde{g}(x) - g(x)| = \left| \sum h_{y_i}(x)(f_{y_i}(x) - g(x)) - h_V(x)f(x) \right| < \varepsilon + h_V(x)f(x) \leq 2\varepsilon$$

Thus $\|\tilde{g} - g\|_\infty \leq 2\varepsilon$. It follows that g can be approximated to arbitrary precision by elements of \mathfrak{a} , so $g \in \mathfrak{a}$, since \mathfrak{a} is closed. We conclude \mathfrak{a} consists of all functions which vanish on C . Thus closed ideals of $C(K)$ are in one-to-one correspondence with closed sets of K . But this means that the maximal ideals are precisely in one to one correspondence with the points of K . Thus the points in the character space of $C(K)$ is in one-to-one correspondence with the points in K .

Given a unital Banach algebra A , and $M \in A$, by duality we obtain a function $\widehat{M} : \Phi_A \rightarrow \mathbb{C}$, given by $\widehat{M}(\phi) = \phi(M)$. It is natural to equip Φ_A with the weakest topology such that each of the functions \widehat{M} is continuous. Thus we obtain an algebra homomorphism $\Gamma : A \rightarrow C(\Phi_A)$, given by $\Gamma(M) = \widehat{M}$.

Theorem 9.6. *For any unital Banach algebra A , Φ_A is a compact, Hausdorff space.*

Proof. Let us first verify the Hausdorff condition. Given $\phi, \psi \in \Phi_A$, there is $M \in A$ such that $\phi(M) \neq \psi(M)$. If we pick disjoint neighborhoods U_0 and V_0 of $\phi(x)$ and $\psi(x)$, then

$$U = \widehat{M}^{-1}(U_0) \quad \text{and} \quad V = \widehat{M}^{-1}(V_0)$$

are disjoint open sets in Φ_A separating ϕ and ψ . Thus Φ_A is a Hausdorff space.

Compactness is a little trickier. The space Φ_A can be viewed as a subset of A^* , since characters are linear maps, and the topology on Φ_A is the relative topology on A^* when A^* is equipped with the weak $*$ topology. Since $\|\phi\| \leq 1$ for each $\phi \in \Phi_A$, Φ_A is a bounded subset of A^* . If we can show it is also weak $*$ closed, then the Banach-Alaoglu theorem will imply that Φ_A is compact.

To show Φ_A is weak- $*$ closed, it suffices to show any functional α in the weak $*$ closure of Φ_A is a character. Fix $M, N \in A$, and $\varepsilon > 0$. Consider

$$W = \{\psi \in A^* : |\psi - \alpha|(z) < \varepsilon \text{ for } z \in \{1, M, N, MN\}\}$$

Then W is a weak- $*$ neighbourhood of α , and thus contains some $\phi \in \Phi_A$. Thus

$$|\phi(1) - \alpha(1)| = |1 - \alpha(1)| < \varepsilon.$$

Thus we find that $\alpha(1) = 1$, since ε was arbitrary. Furthermore,

$$\begin{aligned} \alpha(MN) - \alpha(M)\alpha(N) &= [\alpha(MN) - \phi(MN)] + [\phi(M)\phi(N) - \alpha(M)\alpha(N)] \\ &= [\alpha(MN) - \phi(MN)] + [\phi(N) - \alpha(N)]\phi(M) \\ &\quad + [\phi(M) - \alpha(M)]\alpha(N) \end{aligned}$$

Hence, since ϕ is a complex homomorphism,

$$|\alpha(MN) - \alpha(M)\alpha(N)| < (1 + |\phi(M)| + |\alpha(N)|)\varepsilon \leq (1 + \|M\| + \|\alpha\|\|N\|)\varepsilon$$

Letting $\varepsilon \rightarrow 0$, we find $\alpha(MN) = \alpha(M)\alpha(N)$. Thus α is a character, so Φ_A is closed, and thus compact. \square

Theorem 9.7. *If $f : A \rightarrow B$ is an algebra homomorphism, then the map $f^* : \Phi_B \rightarrow \Phi_A$ defined such that for any $M \in A$,*

$$f^*(\phi)(M) = \phi(f(M))$$

is continuous.

Proof. The open sets in Φ_A are generated for $M \in A$, $\lambda \in \mathbf{C}$, and $\varepsilon > 0$, by sets of the form

$$U = \{\phi \in \Phi_A : |\phi(M) - \lambda| < \varepsilon\} = \widehat{M}^{-1}(B_\varepsilon(\lambda)).$$

The inverse image of U under f^* is precisely

$$\{\psi \in \Phi_B : |(\psi \circ f)(M) - \lambda| < \varepsilon\} = \widehat{f(M)}^{-1}(B_\varepsilon(\lambda)).$$

The map is therefore continuous. □

Let us now list the main properties of the Gelfand transform, which make it so useful.

Theorem 9.8. *Let A be a unital Banach algebra. Then for any $M \in A$,*

$$\widehat{M}(\Phi_A) \subset \sigma(M).$$

If A is commutative, then

$$\widehat{M}(\Phi_A) = \sigma(M).$$

Proof. If $\phi : A \rightarrow \mathbf{C}$, and $\phi(M) = \lambda$, then $\lambda - M$ lies in some maximal ideal of A , and thus cannot be invertible, so $\lambda \in \sigma(M)$. Conversely, if A is commutative, and $\lambda \in \sigma(M)$, then $\lambda - M$ lies in some maximal ideal of A , and so there is some character $\phi : A \rightarrow \mathbf{C}$ with $\phi(\lambda - M) = 0$, i.e. $\phi(M) = \lambda$. □

Example. *The Gelfand transform can be quite useless when studying a non-commutative algebra A . For instance, the algebra $M_n(\mathbf{C})$ of $n \times n$ matrices has only a single ideal, namely the set $\{0\}$. To see this, suppose \mathfrak{a} is an ideal of $M_n(\mathbf{C})$ containing some non-zero matrix M . It follows from the elementary theory of matrices that there are invertible matrices N and L such that $NML \in \mathfrak{a}$ is a diagonal matrix, non-zero somewhere on the diagonal. Multiplying on the left by an appropriate matrix, we may assume that there is only*

a single entry on the diagonal. Multiplying on the left by permutation matrices and then considering linear combinations of what we get then shows that \mathfrak{a} contains all diagonal matrices in $M_n(\mathbb{C})$. But as we have seen above, this implies that \mathfrak{a} contains all matrices in $M_n(\mathbb{C})$. Thus $M_n(\mathbb{C})$ is a simple algebra, it contains no non-trivial ideals. Consequently, we find that every non-zero algebra homomorphism $F : M_n(\mathbb{C}) \rightarrow A$ must be injective, for the kernel is a proper ideal for $M_n(\mathbb{C})$. Thus the character space of $M_n(\mathbb{C})$, for $n > 1$, must be \emptyset , because there is no injective algebra homomorphism from $M_n(\mathbb{C})$ to \mathbb{C} (there is not even an injective linear map).

Since $\widehat{M}(\Phi_A) \subset \sigma(M)$ for any M in a unital Banach algebra A , it follows that

$$\|\Gamma(M)\|_{L^\infty(\Phi_A)} = \rho(M) \leq \|M\|.$$

Thus Γ is a contractive algebra homomorphism from A into $C(\Phi_A)$, whose kernel is the Jacobson radical of A , the intersection of all its maximal ideals. Thus Γ is injective if and only if A is semisimple.

Example. An example of a commutative Banach algebra which is not semisimple is the ring A of all matrices of the form

$$\begin{pmatrix} a & b \\ 0 & a \end{pmatrix}.$$

Any matrix of the form above, where $a \neq 0$, is invertible in A , and the complement of this set is an ideal. Thus this ring is local. The Gelfand space Φ_A in this case is just a single point, and a matrix of the form above is mapped to the function taking the value a at this point.

For general Banach $*$ algebras, Gelfand theory does not necessarily react well with involution. We would want to have the property that

$$\widehat{M^*} = \overline{\widehat{M}},$$

for all $M \in A$. This condition is equivalent to $\Gamma : A \rightarrow C(\Phi_A)$ being a homomorphism of Banach $*$ algebras. A Banach $*$ algebra with this property is called *symmetric*. An example consequence is that for any $M \in A$,

$$\widehat{M^*M} = |\widehat{M}|^2,$$

and if M is unitary, then $|\widehat{M}| = 1$.

Lemma 9.9. *Any unital C^* algebra A is symmetric.*

Proof. We have already seen that $\sigma(M) \subset \mathbf{R}^+$ for any self-adjoint $M \in A$, so the theorem follows immediately from the holomorphic functional calculus and the fact that $\widehat{M}(\Phi_A) \subset \sigma(M)$. Alternatively, we can give a more elementary proof. Consider a self-adjoint element $M \in A$. If $\phi : A \rightarrow \mathbf{C}$ is a character, and $\phi(M) = a + ib$, then $\phi(M + it) = a + i(b + t)$. But we have

$$\begin{aligned} |\phi(M + it)|^2 &= a^2 + (b + t)^2 \\ &\leq \|M + it\|^2 \\ &= \|(M + it)(M - it)\| \\ &= \|M^2 + t^2\| \leq \|M\|^2 + t^2. \end{aligned}$$

Subtracting t^2 from each side of the equation gives a uniform bound on $a^2 + b^2 + 2bt$ for all t , which is only possible if $b = 0$. \square

Remark. The algebra $A(\mathbf{D})$ with involution $f^*(z) = \overline{f(\bar{z})}$, is not symmetric. Neither is $L^1(\mathbf{Z})$, where $(a^*)_n = \bar{a}_n$. But if we instead consider the convolution $(a^*)_n = \overline{a_{-n}}$, then $L^1(\mathbf{Z})$ is a symmetric algebra, which is one of the reasons why this involution is so much more useful than pointwise conjugation.

Theorem 9.10. *If A is a commutative, symmetric Banach $*$ algebra, then the image of the Gelfand transform is a dense subset of $C(\Phi_A)$.*

Proof. Let $B = \Gamma(A)$. Then B is closed under conjugation because A is symmetric. Since characters are determined by their values on elements of A , elements of B separate points in Φ_A . It also contains constants since $\hat{1} = 1$. Thus by the Stone-Weirstrass theorem, B is dense in $C(\Phi_A)$. \square

If A is a commutative unital Banach algebra, the spectral radius formula implies that Γ is an isometry if and only if for any $M \in A$,

$$\|M^2\| = \|M\|^2.$$

This is true, in particular, if A is a commutative C^* algebra. To see this, using the fact that A is symmetric, we calculate that

$$\|M\|^2 = \|M^*M\| = \|\widehat{M^*M}\|_{L^\infty(\Phi_A)} = \|\widehat{M}^*\widehat{M}\|_{L^\infty(\Phi_A)} = \|(\widehat{M})^2\|_{L^\infty(\Phi_A)} = \|M^2\|.$$

Thus *every commutative, unital C^* algebra is semisimple*. Moreover, the Gelfand transform $\Gamma : A \rightarrow C(\Phi_A)$ is then an *isometric $*$ -isomorphism* between A and $C(\Phi_A)$. Thus *every commutative C^* algebra is naturally isomorphic to $C(K)$ for some compact space K* . This is the commutative case of the Gelfand-Naimark theorem.

The commutative Gelfand Naimark theorem has two important applications. The first is a proof of the existence of the *Stone-Cech compactification* of a completely regular topological space. Recall that a space X is *completely regular* if it is T_1 , and if for each closed C and $x \in X - C$, there is $f \in C(X)$ which satisfies $f(x) = 1$, but vanishes on C . Urysohn's lemma tells us all normal spaces are completely regular.

Theorem 9.11. *Let X be a completely regular space. Then there is a compact Hausdorff space βX with an embedding $i : X \rightarrow \beta X$ onto a dense subset of βX satisfying the following universal property: For any continuous map $j : X \rightarrow K$ into a compact Hausdorff space, there is a unique $F : \beta X \rightarrow K$ such that $j = F \circ i$. The universal property implies βX is unique up to a homeomorphism, and it is called the Stone-Cech compactification of X .*

Proof. The space $C_b(X)$ is a commutative unital C^* algebra, so the Gelfand transform gives an isomorphism of $C_b(X)$ with $C(\beta X)$, where $\beta X = \Phi_{C_b(X)}$ is the Character space of $C_b(X)$. For each $x \in X$, the map $\phi_x(f) = f(x)$ is a character on $C_b(X)$, and thus corresponds to a unique point $i(x) \in \beta X$. We thus obtain an injective map $i : X \rightarrow \beta X$. We claim it is also continuous. Indeed, the open sets of βX are generated, for each $f \in C_b(X)$, $\lambda \in \mathbb{C}$, and $\varepsilon > 0$, by sets of the form

$$U = \{x \in \beta X : |\hat{f}(x) - \lambda| < \varepsilon\}.$$

The inverse image of U by i is precisely

$$\{x \in X : |f(x) - \lambda| < \varepsilon\},$$

which is certainly open. Thus i is continuous. We will now show $i : X \rightarrow \beta X$ satisfies the required universal properties.

To begin with, note that $\hat{f} \circ i = f$ for any $f \in C_b(X)$. Since $C_b(X)$ is a C^* algebra, $\Gamma : C_b(X) \rightarrow C(\beta X)$ is an isomorphism, so every continuous, scalar-valued function on βX is the extension of a unique bounded, continuous function on X with the same L^∞ norm. It follows that $i(X)$ is dense in $\beta(X)$, e.g. by Urysohn's lemma.

Next, we show $i : X \rightarrow \beta X$ is an embedding. Fix a net $\{x_\alpha\}$ in X , let $x \in X$, and suppose $\{x_\alpha\}$ does not converge to x . Thinning this net if necessary, we may find an open set U containing x , but no elements of the net $\{x_\alpha\}$. Since X is completely regular, we may find $f \in C_b(X)$ which vanishes outside of U , but with $f(x) = 1$. Then $\hat{f}(i(x_\alpha)) = 0$, but $\hat{f}(i(x)) = 1$. Thus $\{i(x_\alpha)\}$ cannot converge to $i(x)$. Thus we have shown i is an embedding.

Finally, let K be any compact Hausdorff space, and let $j : X \rightarrow K$ be a continuous map. Then j induces an algebra homomorphism $j^* : C(K) \rightarrow C_b(X)$, and thus a continuous map from the character space of $C_b(X)$ to the character space of $C(K)$, i.e. a continuous map $F : \beta X \rightarrow K$ such that if $x \in \beta X$ corresponds to a character $\phi : C_b(X) \rightarrow \mathbb{C}$, then for any $g \in C(K)$,

$$g(F(x)) = \phi(j^*(g)).$$

In particular, if $x \in X$, then

$$g(F(i(x))) = \phi_x(j^*(g)) = j^*(g)(x) = g(j(x)).$$

Since g was arbitrary, it follows that $F \circ i = j$. That F is the unique continuous map making the diagram commute follows because $i(X)$ is dense in βX . Thus we have verified the universal property. \square

For normal elements of a C^* algebra, we now use the Gelfand-Naimark theorem to extend the holomorphic functional calculus to a *continuous functional calculus*. To begin with, note that A is a commutative, unital Banach $*$ algebra containing an element M such that one of three things is true:

- A is generated as a unital algebra by M .
- M is invertible, and A is generated as an algebra by M and M^{-1} .
- A is a symmetric algebra, and A is generated as a unital algebra by M and M^* .

Then $\widehat{M} : \Phi_A \rightarrow \sigma(M)$ is a homeomorphism. Since Φ_A is compact, it suffices to show that \widehat{M} is injective. But this follows because, in each of the three cases above, if two characters agree at M , they agree on the entirety of M .

Example. Consider $\delta^1 \in L^1(\mathbf{Z})$. Then $\sigma(\delta^1) = \mathbf{T}$. We have $\sigma(\delta^1) \subset \mathbf{D}$ because $\|\delta^1\|_{L^1(\mathbf{Z})} = 1$. The function δ^1 has an inverse δ^{-1} , and by the spectral theorem and the fact that $\|\delta^{-1}\|_{L^1(\mathbf{Z})} = 1$, $\sigma(\delta^1)^{-1} = \sigma(\delta^{-1}) \subset \mathbf{D}$, so we conclude $\sigma(\delta^1) \subset \mathbf{T}$. In fact, $\sigma(\delta^1)$ is equal to \mathbf{T} . To see this, let us try and invert $\lambda\delta^0 - \delta^1$, for some $|\lambda| = 1$. If $c * (\lambda\delta^0 - \delta^1) = \delta^0$, then

$$\lambda c_n - c_{n-1} = \begin{cases} 1 & n = 0 \\ 0 & n \neq 0 \end{cases}$$

Solving these equations recursively, we find $c_n = \lambda^{-n}c_0$ for $n \geq 0$, and $c_{-n} = \lambda^{n-1}c_{-1}$. We must have $c \in L^1(\mathbf{Z})$, so that

$$|c_0| \sum_{k=0}^{\infty} 1 + |c_{-1}| \sum_{k=1}^{\infty} 1 < \infty$$

Hence $c_0 = c_{-1} = 0$, so that $c = 0$, which is impossible if we also want $c * (\lambda\delta^0 - \delta^1) = \delta^0$ to hold. Thus $\sigma(\delta^1) = \mathbf{T}$.

Since δ^1 generates $L^1(\mathbf{Z})$, it follows that the Gelfand space of $L^1(\mathbf{Z})$ is homeomorphic to \mathbf{T} . We have seen that $\Gamma(\sum a_n \delta^n) = \sum a_n z^n$. Taking uniform limits of this result gives that $\Gamma : L^1(\mathbf{Z}) \rightarrow C(\mathbf{T})$ is precisely the inverse Fourier transform on the Torus. Conversely, more sophisticated arguments show that for any locally compact abelian group G , the space $L^1(G)$, equipped with convolution (against the Haar measure) as multiplication has Gelfand space equal to \hat{G} , the dual group of complex-valued characters on G , and $\Gamma : L^1(G) \rightarrow C(\hat{G})$ is precisely the abstract Fourier transform of functions on G . Thus the Gelfand transform is one way to think of the Fourier transform.

Now suppose that A is a commutative unital C^* algebra, either generated by M , or generated by M and M^{-1} for some invertible M (this is obtained in practice, by starting with a C^* algebra, and then considering the smallest C^* algebra generated by a normal element). Then $\widehat{M} : \Phi_A \rightarrow \sigma(M)$ is a homeomorphism. We thus have an isometric $*$ isomorphism

$$A \cong C(\Phi_A) \cong C(\sigma(M)).$$

For $f \in C(\sigma(M))$, we let $f(M) \in A$ denote the image of this map. Thus for any character $\phi : A \rightarrow \mathbf{C}$,

$$\phi(f(M)) = f(\phi(M)).$$

This equation implies that $\sigma(f(M)) = f(\sigma(M))$, so we have a continuous functional calculus. Thus we have a continuous functional calculus on A . It is simple to see that $z(M) = M$, and $1(M) = 1$, and $f \mapsto f(M)$ is the unique $*$ algebra homomorphism with this property, since functions of the form

$$f(z) = \sum c_{\alpha\beta} z^\alpha \bar{z}^\beta$$

are dense in $C(\sigma(M))$, and we see that $f(M) = \sum c_{\alpha\beta} M^\alpha (M^*)^\beta$. In particular, we see that if $f(z) = \sum c_k z^k$, then $f(M) = \sum c_k M^k$, and so this calculus extends the holomorphic functional calculus.

The key to Gelfand theory is noticing that characters of Banach algebras naturally reflect the structure of the Banach algebra in question. Understanding the character space leads to a natural understanding of the invertibility of the elements of the algebra. For instance, here is a result on the invertibility of Fourier series initially due to Wiener, whose proof was incredibly difficult without the technology that Gelfand introduced.

Example. Consider a formal trigonometric series on \mathbf{T} of the form

$$f(z) = \sum a_n z^n$$

where $\sum |a_n| < \infty$. Then $f = \Gamma(a)$, where $a \in L^1(\mathbf{Z})$ is given by the sequence $\{a_n\}$. Now suppose f is non-vanishing on \mathbf{T} . Since $\sigma(a) = f(\mathbf{T})$, this means that a is invertible in $L^1(\mathbf{Z})$. Thus there exists $b \in L^1(\mathbf{Z})$ such that $a * b = \delta$. Thus $\Gamma(b) = 1/f$, i.e.

$$\frac{1}{f(z)} = \sum b_n z^n.$$

Similar results can be applied to non-continuous functions by approximation, and to higher dimensional trigonometric series, in which case the Gelfand space are the higher dimensional torii \mathbf{T}^n .

Here's another example from the complex analysis of several variables.

Example. Consider the commutative algebra $A(\mathbf{D})$. We also contend that $\Phi_{A(\mathbf{D})}$ can be identified with \mathbf{D} . The identity function z has norm 1 in this space, implying that if $\phi : A(\mathbf{D}) \rightarrow \mathbf{C}$ is any character, then $\phi(z) = w \in \mathbf{D}$. Then

$$\phi\left(\sum_{k=0}^N a_k z^k\right) = \sum_{k=0}^N a_k w^k$$

and these polynomials are dense in $A(\mathbf{D})$, from which we conclude $\phi(f) = f(w)$. Thus any character is given by evaluation at a point in \mathbf{D} .

Let us use this calculation to obtain an analytic version of the nullstellensatz in algebraic geometry. Consider n functions $f_1, \dots, f_n \in A(\mathbf{D})$ whose common nullset $Z(f_1, \dots, f_n)$ is empty. Then we claim there are $g_1, \dots, g_n \in A(\mathbf{D})$ such that $f_1 g_1 + \dots + f_n g_n = 1$. It is clear that this is equivalent to the ideal relation

$$(f_1, \dots, f_n) = A(\mathbf{D})$$

If this were not true, then (f_1, \dots, f_n) would be contained in a closed, maximal ideal, and would therefore be annihilated by some $\phi \in \Phi_{A(\mathbf{D})}$. But ϕ corresponds to evaluation at some $w \in \mathbf{D}$, so $f_i(w) = 0$ for all i , which gives a contradiction. One can apply Runge's theorem to classify the character spaces of arbitrary spaces of the form $A(\Omega)$, but we leave this as an exercise to the reader.

If M and N are operators such that $MN = NM$, then certainly $M^*N^* = N^*M^*$, but it is not necessarily true that $MN^* = N^*M$ or $M^*N = NM^*$. For instance, if M is not a normal operator, and we take $M = N$. But if M and N are normal, then the result is true. In fact, much more is true.

Theorem 9.12 (Fuglede-Putnam-Rosenblum). *Let A be a C^* algebra, and consider normal elements $M, N \in A$, and a third element $L \in A$ such that $ML = LM$. Then $M^*L = LN^*$.*

Proof. Without loss of generality, assume A is unital. If $ML = LM$, then for all $k > 0$, $M^k L = L N^k$, and so by taking power series, we conclude that for any $\lambda \in \mathbf{C}$,

$$e^{\lambda M} L = L e^{\lambda N},$$

so that $L = e^{\lambda M} L e^{-\lambda N}$. Thus

$$e^{\lambda M^*} L e^{-\lambda N^*} = e^{\lambda(M^* - M)} L e^{\lambda(N - N^*)}.$$

Note that $e^{\lambda(M^* - M)}$ and $e^{\lambda(N - N^*)}$ are both unitary, and so if we set

$$f(\lambda) = e^{\lambda M^*} L e^{-\lambda N^*} = L + \lambda(M^* L - L N^*) + O(\lambda^2),$$

then we see $f : \mathbf{C} \rightarrow A$ is an entire function of λ , and $\|f\|_{L^\infty(\mathbf{C})} \leq \|L\|$ is uniformly bounded. Thus we conclude f is a constant function. But this means that the first order term of the power series vanishes, i.e. $M^* L = L N^*$. \square

9.2 The Non-Unital Gelfand Space

Let us also address the non-unital case of Gelfand theory. In the Gelfand theory, we have a little trouble defining ideals. We call an ideal \mathfrak{a} (an additive subgroup closed under multiplication) for a non-unital algebra A *modular* if there is $N \in A$ such that $MN - M, NM - M \in \mathfrak{a}$ for all $M \in A$. If \mathfrak{a} is modular, it follows that A/\mathfrak{a} contains an identity. Zorn can show us maximal modular ideals exist in any non-unital algebra. We need to edit the proof which shows a maximal ideal is closed, which relies on an algebraic trick.

Lemma 9.13. *A maximal modular ideal in an algebra A is closed.*

Proof. If a maximal ideal \mathfrak{a} was not closed, then we would have $\bar{\mathfrak{a}} = A$. Let N be a right modular identity for \mathfrak{a} . Then there is $M \in \mathfrak{a}$ with $\|M - N\| < 1$, so

$$\begin{aligned} N &= (N - M) + M = \sum_{k=1}^{\infty} (N - M)^k - \sum_{k=2}^{\infty} (N - M)^k + M \\ &= \sum_{k=1}^{\infty} (N - M)^k - \left[\sum_{k=1}^{\infty} (N - M)^k \right] (N - M) + M \\ &= \underbrace{\left[\left(\sum_{k=1}^{\infty} (N - M)^k \right) M + M \right]}_{\in \mathfrak{a}} - \underbrace{\left[\left(\sum_{k=1}^{\infty} (N - M)^k \right) N - \left(\sum_{k=1}^{\infty} (N - M)^k \right) \right]}_{\in \mathfrak{a}} \end{aligned}$$

Given any $L \in A$, $LM - M \in \mathfrak{a}$, implying $M \in \mathfrak{a}$. This implies $\mathfrak{a} = A$, an impossibility. \square

Lemma 9.14. *The kernel of any nonzero algebra homomorphism from A to \mathbf{C} is a maximal modular ideal, and in a commutative non-unital algebra, any maximal modular ideal is the kernel of some non-zero algebra homomorphism.*

Proof. If $\phi : A \rightarrow \mathbf{C}$ is a non-zero algebra homomorphism, then ϕ is surjective, for if $\phi(M) \neq 0$, then $\phi(\mathbf{C} \cdot x) = \mathbf{C}$. We claim that if $\mathfrak{a} = \ker(\phi)$, and if $\phi(M) = 1$, then M is a modular identity for \mathfrak{a} , because

$$\phi(MN - N) = \phi(M)\phi(N) - \phi(N) = \phi(N) - \phi(N) = 0$$

Implying $MN - N \in \mathfrak{a}$ for each N , so \mathfrak{a} is a modular ideal. The first isomorphism theorem gives us an isomorphism from A/\mathfrak{a} to \mathbf{C} , so \mathfrak{a} is maximal, since A/\mathfrak{a} is a field.

If A is commutative, and \mathfrak{a} is a maximal modular ideal, then A/\mathfrak{a} is an algebra with identity, whose only ideals consists of (0) and A/\mathfrak{a} , since \mathfrak{a} is maximal, so A/\mathfrak{a} is a field, and is therefore isometric to \mathbf{C} , giving us an algebra homomorphism. \square

Given this result, we may therefore consider the Gelfand transform of a non-unital algebra, into the character space of all nonzero non-unital algebra homomorphisms. Unitizing A and considering the Gelfand theory on $A^\#$ only adds a single extra character, which normally does not introduce any complications.

Lemma 9.15. *If A is an algebra without identity, and ϕ is a character on A , then there is a unique non-zero character $\tilde{\phi}$ defined on $A^\#$ which extends ϕ .*

Proof. Define $\tilde{\phi}(M + \lambda) = \phi(M) + \lambda$. Then $\tilde{\phi}$ is linear, $\tilde{\phi}(1) = 1$, and

$$\begin{aligned}\tilde{\phi}(M + \lambda)\tilde{\phi}(N + \gamma) &= [\phi(M) + \lambda][\phi(N) + \gamma] \\ &= \phi(MN) + \phi(\lambda N) + \phi(\gamma M) + \lambda\gamma \\ &= \tilde{\phi}((M + \lambda)(N + \gamma))\end{aligned}$$

The uniqueness of the extension follows from the fact any maximal modular ideal of A can be uniquely extended to a maximal ideal on $A^\#$, for the projection $M + \lambda \mapsto M$ maps ideals to ideals, and the only ideals of \mathbf{C} are (0) and (1) . \square

Corollary 9.16. *We have*

$$\Phi_{A^\#} = \Phi_A \cup \{\phi_\infty\},$$

where

$$\phi_\infty(M + \lambda) = \lambda.$$

We have $\sigma(M) \subset \widehat{M}(\Phi_A) \cup \{0\}$, with equality if A is commutative.

We can equip Φ_A with the weak topology in the same way that we introduce a topology in the unital case. In this situation, however, Φ_A is only *locally compact*, not compact, with $\Phi_{A^\#}$ corresponding to the one point compactification.

Theorem 9.17. *The Gelfand topology on Φ_A is Hausdorff and locally compact, and we have*

$$\Gamma : A \rightarrow C_0(\Phi_A).$$

Proof. This follows from the fact that $\Phi_{A^\#}$ is Hausdorff compact, and Φ_A is obtained from $\Phi_{A^\#}$ by removing a single point. Certainly it follows from this that Φ_A is Hausdorff. Moreover, for any character ϕ , we can find an open set $U \subset \Phi_{A^\#}$ containing ϕ and an open set $V \subset \Phi_{A^\#}$ containing ϕ_∞ such that the closures of U and V are disjoint. But this means that $\overline{U} \subset \Phi_A$ is a compact neighborhood of ϕ . Thus Φ_A is locally compact. \square

Example. *Suppose X is a locally compact, but not compact, space. Then $C_0(X)$ is a non-unital C^* algebra. Its unitization is then a unital C^* algebra, and thus corresponds to $C(X^*)$ for some compact set X^* . We have already seen that X^* is equal to the one point compactification of X .*

As with the holomorphic functional calculus on non unital Banach algebras, given a normal element M of a non-unital C^* algebra A , we can define $f(M)$ for any $f \in C(\sigma(M))$, provided that $f(0) = 0$.

Chapter 10

Positivity

The model example of this chapter is a positive semidefinite operator on a Hilbert space.

Theorem 10.1. *Let H be a Hilbert space, and suppose $T : H \rightarrow H$ is a bounded, self-adjoint operator. Then the following two conditions are equivalent:*

- $\sigma(T) \subset [0, \infty)$.
- For any $x \in H$, $\langle Tx, x \rangle \geq 0$.

If these properties are true, we say T is positive semidefinite.

Proof. Suppose $\sigma(T) \subset [0, \infty)$. If we let $f : \sigma(T) \rightarrow [0, \infty)$ be defined by setting $f(x) = \sqrt{x}$, and we set $S = f(T)$, then S is also self-adjoint, and $T = S^2$. It follows that for any $x \in H$,

$$\langle Tx, x \rangle = \langle S^2x, x \rangle = \langle Sx, Sx \rangle \geq 0.$$

Conversely, suppose that for all $x \in H$,

$$\langle Tx, x \rangle \geq 0.$$

Then for any $\lambda < 0$, we have

$$\langle (T - \lambda)x, x \rangle = \langle Tx, x \rangle - \lambda\|x\|^2 \geq -\lambda\|x\|^2 \gtrsim \|x\|^2.$$

Since $|\langle (T - \lambda)x, x \rangle| \lesssim \|(T - \lambda)x\|\|x\|$, it follows that $\|(T - \lambda)x\| \sim \|x\|$, so $T - \lambda$ is an isomorphism of H . Thus $\sigma(T) \subset [0, \infty)$. \square

Now let A be an arbitrary C^* algebra. A self-adjoint element M is *positive* if $\sigma(M) \subset [0, \infty)$. In shorthand, we write $M \geq 0$. The set of all positive elements in an algebra A is denoted A_+ . We write $M \leq N$ if $N - M \geq 0$. Thus identifying positive elements gives us a partial ordering on the set of self-adjoint elements of A .

Example. Let X be a locally compact space. Then $f \in C_0(X)$ is self-adjoint if and only if f is real-valued, and positive if and only if $f \geq 0$, since $\sigma(f) = \{0\} \cup f(X)$. For a measure space X , $f \in L^\infty(X)$ is self-adjoint if and only if f is almost-everywhere real valued, and positive if and only if $f \geq 0$ almost everywhere.

Proposition 10.2. Let A be a C^* algebra. If $M \in A_+$, there is a unique $N \in A_+$ for which $M = N^2$, which we denote by \sqrt{M} .

Proof. Consider the positive square root map $f(z) = \sqrt{z}$. This is well defined on $\sigma(M)$ since $\sigma(M) \subset [0, \infty)$, and $f \in C(\sigma(M))$, so we may consider $N = f(M)$. Now

$$N^2 = f(M)^2 = (f^2)(M) = (\text{id}_{\sigma(M)})(M) = M.$$

Denote this element by \sqrt{M} . If $N^2 = M$ for any other $N \in A_+$, then

$$NM = N^3 = (N^2)N = MN$$

Thus N commutes with M , and so it follows that N commutes with $f(M)$ for any $f \in C(\sigma(M))$, since $f(M)$ lies in the unital C^* algebra generated by M . Thus the smallest C^* algebra B containing N and $f(M)$ is commutative. Applying the Gelfand transform, since $N \in A_+$, $\widehat{N}(\Phi_B) \geq 0$, and $\widehat{N}^2 = M$. This means that $\widehat{N} = \sqrt{\widehat{M}} = \widehat{f(M)}$. Thus $\widehat{f(M)} = \widehat{N}$, so $f(M) = N$. \square

We denote the unique positive square root of $M \in A_+$ by \sqrt{M} . If self-adjoint operators behave somewhat like real numbers, positive operators behave like positive real numbers.

Proposition 10.3. Let A be a C^* algebra, and let $M \in A$ be self adjoint. Then there are unique $M_+, M_- \in A_+$ for which $M = M_+ - M_-$, and $M_+M_- = 0$.

Proof. Let $f_+(t) = \max(t, 0)$, and $f_-(t) = -\min(t, 0)$. Then f_+ and f_- are continuous maps, and we may consider $f_+(M)$, $f_-(M)$. Since

$$f_+ - f_- = \text{id}_{\sigma(M)}$$

$f_+(M) - f_-(M) = M$. By the spectral mapping theorem,

$$\sigma(f_+(M)) = f_+(\sigma(M)) \subset [0, \infty) \quad \sigma(f_-(M)) = f_-(\sigma(M)) \subset [0, \infty)$$

so $f_+(M), f_-(M) \geq 0$. Since $f_+f_- = 0$, $f_+(M)f_-(M) = 0$. Uniqueness follows by a Gelfand transform calculation analogous to the proof of the uniqueness of square roots. \square

Lemma 10.4. *Let A be a C^* algebra. If M is self-adjoint, $t \geq 0$, and $\|M - t\| \leq t$, then M is positive. If M is positive, and $\|M\| \leq t$, then $\|M - t\| \leq t$.*

Proof. Suppose $\|M - t\| \leq t$. Then $\sigma(M - t) \in [-t, t]$, so

$$\sigma(M) = \sigma(M - t) + t \in [0, 2t]$$

Thus $M \in A_+$. Conversely, if $M \in A_+$ and $\|M\| \leq t$, then $\sigma(M) \subset [0, t]$, and

$$\sigma(M - t) = \sigma(M) - t \subset [-t, 0]$$

which means $\|M - t\| = r(M - t) \leq t$. \square

The next result shows that A_+ is a cone in A .

Proposition 10.5. *If A is a C^* algebra, then*

- (a) A_+ is closed.
- (b) If $M, N \in A_+$, then $M + N \in A_+$.
- (c) If $M \in A_+$, and $t \geq 0$, then $tM \in A_+$.
- (d) $A_+ \cap (-A_+) = (0)$.

Proof. Suppose that $N \in \overline{A_+}$, choose $K > \|N\|$. Fix $\varepsilon > 0$ and pick $M \in A_+$ with $\|N - M\| < \varepsilon$. Then if ε is chosen small enough, then $\|M\| < K$, and

$$\|N - K\| = \|N - M\| + \|M - K\| \leq \varepsilon + K$$

Letting $\varepsilon \rightarrow 0$, we find $N \in A_+$. To prove (b), we apply the inequality in the lemma to conclude

$$\|M + N - (\|M\| + \|N\|)\| \leq \|M - \|M\|\| + \|N - \|N\|\| \leq \|M\| + \|N\|$$

So $M + N \in A_+$. Since $\sigma(\lambda M) = \lambda\sigma(M)$, we obtain (c). To prove (d), suppose $M \in A_+ \cap -A_+$. Then $\sigma(M) = \{0\}$, and since M is normal, $\|M\| = r(M) = 0$. \square

The next corollary is an analogy of the result over the real numbers that if $|x| \leq 0$, then $x = 0$.

Lemma 10.6. *If $-M^*M \in A_+$, then $M = 0$.*

Proof. Since

$$\sigma(MM^*) \cup \{0\} = \sigma(M^*M) \cup \{0\}$$

which can be proved by some algebraic tricks, $-MM^* \in A_+$ as well. If we write $M = T + iS$, where T and S are normal, then

$$M^*M + MM^* = (T - iS)(T + iS) + (T + iS)(T - iS) = 2(T^2 + S^2) \in A_+$$

Thus

$$M^*M = 2(T^2 + S^2) - MM^* \in A_+$$

which implies $M^*M = 0$. □

Proposition 10.7. *The following are equivalent*

- (a) $M \in A_+$.
- (b) There is $N \in A_+$ such that $M = N^2$.
- (c) There is $N \in A$ such that $N^*N = M$.

Proof. We have already show (a) implies (b). The proof of (c) from (b) is trivial. To prove (a) from (c), note that if $M = N^*N$, then M is certainly self-adjoint, so we may write $M = M_+ - M_-$. Let $L = NM_-$. Then

$$-L^*L = -M_-N^*NM_- = -M_-MM_- = M_-^3 \in A_+$$

Implying $NM_- = 0$, hence

$$0 = N^*L = N^*NM_- = -M_-^2 \in A_+$$

This implies $M_- = 0$, so $M = M_+ \in A_+$. □

For M lying in a C^* algebra A , we are inclined to define $|M| = \sqrt{MM^*}$, for any $M \in A$, which is a positive element.

Proposition 10.8. *If $M \leq N$, then*

- (a) $M + L \leq N + L$.

$$(b) L^*ML \leq L^*NL.$$

If in addition, $0 \leq M \leq N$, then

$$(c) \|M\| \leq \|N\|.$$

$$(d) \sqrt{M} \leq \sqrt{N}.$$

and if $M, N \in GL(A)$, then

$$(e) 0 \leq N^{-1} \leq M^{-1}$$

Proof. (a) is trivial. To prove (b), note that

$$L^*(N - M)L = L^*\sqrt{N - M}\sqrt{N - M}L = (\sqrt{N - M}L)^*(\sqrt{N - M}L) \in A_+$$

Let us prove (c). If N is positive, then $N \leq \|N\|$, which follows by Gelfand theory. Thus if M and N are positive, then

$$0 \leq M \leq N \leq \|N\|$$

But then $\phi(M) \leq \|N\|$ for each ϕ , so $\|M\| \leq \|N\|$.

Fix $\varepsilon > 0$. Write

$$(\varepsilon + \sqrt{N} + \sqrt{M})(\varepsilon + \sqrt{N} - \sqrt{M}) = T$$

where $T, U \in A$ are self adjoint. By calculation

$$T = \varepsilon^2 + 2\varepsilon\sqrt{N} + N - M \geq \varepsilon^2$$

Thus T is positive and invertible, so $\varepsilon + \sqrt{N} - \sqrt{M}$ is left invertible and therefore invertible. Thus $\varepsilon \notin \sigma(\sqrt{M} - \sqrt{N})$, and thus

$$\sigma(\sqrt{M} - \sqrt{N}) \subset (-\infty, 0)$$

But this implies $\sqrt{N} - \sqrt{M}$ is positive.

If $M \geq \varepsilon$, then $M^{-1} \leq 1/\varepsilon$, which follows by Gelfand theory, since $\phi(M^{-1}) = \phi(M)^{-1}$. Since

$$1 = \sqrt{M}^{-1}M\sqrt{M}^{-1} \leq \sqrt{M}^{-1}N\sqrt{M}^{-1}$$

this yields

$$\sqrt{M}N^{-1}\sqrt{M} \leq 1$$

and by conjugation again,

$$N^{-1} = \sqrt{M}^{-1}\sqrt{M}N^{-1}\sqrt{M}\sqrt{M}^{-1} \leq M^{-1}$$

□

10.1 Positive Approximate Units

We now show that all C^* algebras have *approximate units*, which are bounded approximate identities which are also increasing nets, in the sense of the positive ordering on the C^* algebra.

Lemma 10.9. *If $0 \leq M \leq N$, then $M(1 + M)^{-1} \leq N(1 + N)^{-1}$.*

Proof. Note that

$$M(1 + M)^{-1} = 1 - (1 + M)^{-1} \quad N(1 + N)^{-1} = 1 - (1 + N)^{-1}$$

We know $1 + M \leq 1 + N$, so $(1 + N)^{-1} \leq (1 + M)^{-1}$, so

$$N(1 + N)^{-1} = 1 - (1 + N)^{-1} \leq 1 - (1 + M)^{-1} = M(1 + M)^{-1}$$

which is exactly the inequality we needed. \square

Proposition 10.10. *For a C^* algebra A , the net*

$$\{M \in A_+ : \|M\| < 1\}$$

ordered by the positivity of the algebra, is an approximate unit.

Proof. The set is a net, because any $M, N \in A_+$ are less than or equal to $\max(\|M\|, \|N\|)$. Fix a self adjoint $M \in A$, and let $\varepsilon > 0$. Let B be the (non-unital) C^* algebra generated by M . Then B is commutative, and $\widetilde{M} \in C_0(\Phi_B)$. Let $g \in C_0(\Phi_B)$ be a function such that $0 \leq g < 1$, such that $\|\widetilde{M} - g\widetilde{M}\| < \varepsilon^2$, and set N to be the unique element of B such that $\widetilde{N} = g$. Then $\|N\| \leq 1$, N is positive, and $\|M - MN\| < \varepsilon^2$. Let $N \leq L < 1$. Then, since $\sigma(\sqrt{1 - L}) \in [0, 1]$, $\|\sqrt{1 - L}\| < 1$, and

$$\begin{aligned} \|M - LM\|^2 &= \|\sqrt{1 - L}^2 M\|^2 \leq \|\sqrt{1 - L} M\|^2 = \|M(1 - L)M\| \\ &\leq \|M(1 - N)M\| \leq \|M\|(1 - N)\|M\| < \|M\|\varepsilon^2 \end{aligned}$$

Decreasing ε , we find $\lim LM = M$. Similar results show $\lim ML = M$. \square

If A is an abelian C^* algebra, then the family of self adjoint elements in A forms a real Riesz lattice, which can be easily seen from an isometric representation $A \cong C(X)$. Thus if $0 \leq M \leq N + L$, where $N, L \in A_+$, then there are N_0, L_0 such that $N_0 \leq N$, $L_0 \leq L$, and such that $M = N_0 + L_0$. This theorem does not hold in general nonabelian C^* algebras, but there is a weaker decomposition theorem which does hold.

Proposition 10.11. *If A is a C^* algebra, and*

$$\sum_{i=0}^m M_i^* M_i = \sum_{j=0}^n N_j^* N_j$$

then there are $L_{i,j}$ such that

$$M_i^* M_i = \sum_j L_{i,j}^* L_{i,j} \quad N_j^* N_j = \sum_i L_{i,j}^* L_{i,j}$$

Proof. We may assume A is unital. Let $X = \sum M_i^* M_i = \sum N_j^* N_j$. Write

$$L_{i,j,t} = s$$

and so on and so forth FINISH THIS LATER. □

10.2 Ideals and Quotients of C^* Algebras

10.2.1 C Star Operations on Quotients

If we quotient a C^* algebra by a closed ideal, we certainly get a Banach algebra back. But does this algebra still have an involution, or do we need to apply additional structure to our ideals? Such discussions naturally lead to the Gelfand Naimark construction, establishing that every C^* algebra is isometric to some subalgebra of bounded operators on a Hilbert space. Thus such an endeavor is very fruitful to discuss.

Proposition 10.12. *If \mathfrak{a} is a closed left/right ideal in a C^* algebra A containing a self-adjoint element M , and if $f \in C(\sigma(M))$ satisfies $f(0) = 0$, then $f(M) \in \mathfrak{a}$.*

Proof. Choose a sequence P_i of polynomials which uniformly approximating f , with $P_i(0) = 0$ for all i . If we write

$$P_i = \sum a_{ij} X^i \bar{X}^j$$

then $P_i(M) = \sum a_{ij} X^{i+j} \in \mathfrak{a}$, and since \mathfrak{a} is closed, we find $f(M) \in \mathfrak{a}$. □

Corollary 10.13. *If $M \in \mathfrak{a}$ is self adjoint, then M_+ , M_- , $|M|$, and \sqrt{M} are all in \mathfrak{a} .*

If an ideal is closed under involution, then it certainly must be a two sided ideal, for involution reverses multiplication. We shall show that all two-sided ideals are closed under involution.

Proposition 10.14. *A closed ideal in a C^* -algebra is closed under involution.*

Proof. Given a closed ideal \mathfrak{a} of a C^* algebra A , \mathfrak{a} is a (non-unital) C^* sub-algebra of A . Let $\{E_\alpha\}$ be a two-sided approximate identity for \mathfrak{a} . Then, if $M \in \mathfrak{a}$ is arbitrary,

$$\lim M^* E_\alpha = \lim (E_\alpha M)^* = (\lim E_\alpha M)^* = M^*$$

so M^* is the limit of elements of \mathfrak{a} , and is thus in \mathfrak{a} . \square

Thus if \mathfrak{a} is a closed ideal in a C^* -algebra A , then A/\mathfrak{a} is verified to be a Banach $*$ algebra. It is a little bit more technical to verify that A/\mathfrak{a} is a C^* algebra, but we shall now carry out the details.

Lemma 10.15. *If \mathfrak{a} is a closed ideal of A , and $\{E_\alpha\}$ is the positive BAI for \mathfrak{a} , then A/\mathfrak{a} , then for any $M \in A$,*

$$\|M + \mathfrak{a}\| = \lim \|M - E_\alpha M\| = \lim \|M - ME_\alpha\|$$

Proof. Fix $M \in A$, and let $\{E_\alpha\}$ be the positive BAI for \mathfrak{a} . Choose $\varepsilon > 0$, and let $N \in \mathfrak{a}$ satisfy $\|M + N\| \leq \|M + \mathfrak{a}\| + \varepsilon$. Then

$$\begin{aligned} \limsup_\alpha \|M - E_\alpha M\| &\leq \limsup_\alpha \|(1 - E_\alpha)(M + N)\| + \|E_\alpha N - N\| \\ &\leq \limsup_\alpha \|(1 - E_\alpha)(M + N)\| + \limsup_\alpha \|N - E_\alpha N\| \\ &= \limsup_\alpha \|(1 - E_\alpha)(M + N)\| \\ &\leq \|M + N\| \leq \|M + \mathfrak{a}\| + \varepsilon \end{aligned}$$

and ε was arbitrary, so

$$\limsup_\alpha \|M - E_\alpha M\| \leq \|M + \mathfrak{a}\|$$

But $E_\alpha M \in \mathfrak{a}$, so for each α ,

$$\|M - E_\alpha M\| \geq \|M + \mathfrak{a}\|$$

and since we have bounded the limsup above and below by the same value, it is in fact a limit, and has value $\|M + \mathfrak{a}\|$. Similar results hold for $\|M - ME_\alpha\|$. \square

The lemma makes sense, for E_α approximates elements of \mathfrak{a} as best as possible, so subtracting $E_\alpha M$ subtracts the best approximation of M in \mathfrak{a} , thus giving us the quotient norm.

Theorem 10.16. *If \mathfrak{a} is a closed ideal of a C^* algebra A , then A/\mathfrak{a} is a C^* algebra.*

Proof. Letting $\{E_\alpha\}$ be the positive approximate identity for \mathfrak{a} , we find

$$\begin{aligned}\|M + \mathfrak{a}\|^2 &= \lim_{\alpha} \|M(1 - E_\alpha)\|^2 \\ &= \lim_{\alpha} \|(1 - E_\alpha)M^*M(1 - E_\alpha)\| \\ &\leq \lim_{\alpha} \|(1 - E_\alpha)M^*M\| \\ &= \|M^*M + \mathfrak{a}\| \\ &\leq \|M^* + \mathfrak{a}\| \|M + \mathfrak{a}\|\end{aligned}$$

And it is easy to see that $\|M^* + \mathfrak{a}\| = \|M + \mathfrak{a}\|$, for \mathfrak{a} is closed under involution, so we find

$$\|M + \mathfrak{a}\|^2 = \|M^*M + \mathfrak{a}\|$$

and this is exactly the C^* identity for the quotient algebra. \square

It is often much easier to calculate the norm on A/\mathfrak{a} by applying the theory we've created. If $\pi : A \rightarrow B$ is a C^* homomorphism with kernel \mathfrak{a} , then we obtain an induced diagram of C^* morphisms

$$\begin{array}{ccc} A & \xrightarrow{\pi} & B \\ \downarrow & \nearrow \tilde{\pi} & \\ A/\mathfrak{a} & & \end{array}$$

and we know that $\tilde{\pi}$ is an isometry, since it is injective. Thus

$$\|M + \mathfrak{a}\| = \|\pi(M)\|$$

and this is normally easier to calculate than either of the two equations we can use to calculate the quotient norm.

Example. *We know the closed ideals of $C(X)$ are of the form*

$$\mathfrak{a}_K = \{f \in C(X) : (\forall x \in K : f(x) = 0)\}$$

for some closed set K . Then $C(X)/\mathfrak{a}_K$ is isometric to $C(K)$, because we have the C^ homomorphism $f \mapsto f|_K$, and \mathfrak{a}_K is the kernel.*

Example. Let H be a Hilbert space, and \mathfrak{a} a closed ideal of $B(H)$. If $\mathfrak{a} \neq (0)$, we claim $K(H) \subset \mathfrak{a}$. For any $x, y \in H$, let $x \otimes y \in B(H)$ be the map

$$z \mapsto \langle z, y \rangle x$$

It is clear that $K(H)$ is the closed linear hull of $\{x \otimes y : x, y \in H\}$. If $T \in \mathfrak{a}$ and $T \neq 0$, then there are v, w such that $\langle Tv, w \rangle = 1$. Then

$$(x \otimes y) = (x \otimes w) \circ (Tv \otimes y) = (x \otimes w) \circ T \circ (v \otimes y) \in \mathfrak{a}$$

from which we obtain, since $x \otimes y$ was arbitrary, that $K(H) \subset \mathfrak{a}$.

Example. In an infinite dimensional, separable Hilbert space H the only closed ideals of $B(H)$ are (0) , $K(H)$, and $B(H)$. Let \mathfrak{a} be an ideal properly containing $K(H)$. Let $T \in \mathfrak{a} - K(H)$, which is without loss of generality positive. FINISH THIS LATER WHEN BETTER AT SPECTRAL MEASURES.

Thus in an infinite dimensional, separable Hilbert space H , the Calkin Algebra $C(H) = B(H)/K(H)$ has no nontrivial closed ideals.

10.2.2 * Homomorphisms

The natural morphism between C^* algebras is a $*$ homomorphism, an algebra homomorphism T which satisfies $T(M^*) = T(M)^*$.

Proposition 10.17. If $T : A \rightarrow B$ is a $*$ morphism, then $\|T(M)\| \leq \|M\|$ for all M .

Proof. Assume first that A is unital. For any M , $\sigma(T(M)) \subset \sigma(M)$. But then if M is normal, then $T(M)$ is normal, and

$$\|T(M)\| = r(T(M)) \leq r(M) = \|M\|$$

In general, for any M ,

$$\|T(M)\|^2 = \|T(M)T(M)^*\| = \|T(MM^*)\| \leq \|MM^*\| = \|M\|^2$$

so $\|T\| \leq 1$. In the non-unital case, T induces a $*$ -extension between $A^\#$ and $B^\#$, in which we can apply the previous case. \square

Thus every $*$ -morphism is continuous. We shall find that injective $*$ -morphisms are isometries. To prove this, we must first verify a corresponding theorem about continuous functions on a locally compact space, which says that the uniform norm on $C(X)$ is the sharpest algebra norm we can have.

Lemma 10.18. *If $\pi : A \rightarrow B$ is an injective $*$ -morphism between abelian C^* algebras, then the dual map π^* is surjective.*

Proof. Let

$$K = \{\phi|_A : \phi \in \Phi_B\}$$

If $K \neq \Phi_A$, then using the continuous functional calculus, we may find non-zero f and g which vanish at K , for which $f = fg$. Then $1 \in \sigma(g)$, so there is $\phi \in \Phi_B$ such that $\phi(g) = 1$. But this is clearly impossible. \square

Lemma 10.19. *Let $\|\cdot\|$ be a submultiplicative norm on $C(X)$. Then $\|\cdot\|_\infty \leq \|\cdot\|$.*

Proof. Let A be the completion of $C(X)$ with respect to the $\|\cdot\|$ norm. Then A is a Banach algebra, and we obtain

$$\|f\|_\infty = r(f) \leq \|f\|$$

The fact that algebraic facts can hide within them topological facts is essential to the proof. \square

Since arbitrary C^* algebras contain abelian subalgebra on mass, we can almost effortlessly harvest facts about $*$ morphisms on arbitrary algebras from the above lemma.

Proposition 10.20. *An injective, continuous morphism $\pi : A \rightarrow B$ from a C^* algebra to a Banach algebra satisfies*

$$\|M\| \leq \|\pi\| \|\pi(M)\|$$

Proof. Given M , let $A_0 = C^*(M^*M)$. Define a norm on A_0 by letting

$$\|N\| = \|\pi(N)\|$$

Then $\|\cdot\| \leq \|\cdot\|$, and

$$\|M\|^2 = \|M^*M\| \leq \|M^*M\| = \|\pi(M^*M)\| \leq \|\pi\| \|M\| \|\pi(M)\|$$

which yields the claim. \square

Corollary 10.21. *Let $\pi : A \rightarrow B$ be a $*$ homomorphism between C^* algebras. Then $\pi(A)$ is a closed C^* algebra, and π is an isometry if it is injective.*

Proof. If π is injective, it certainly has closed range by the last proposition. But in general, letting $\mathfrak{a} = \ker(\pi)$, π induces an injective map

$$\begin{array}{ccc} A & \xrightarrow{\pi} & B \\ \downarrow & \nearrow \tilde{\pi} & \\ A/\mathfrak{a} & & \end{array}$$

and $\pi(A) = \tilde{\pi}(A)$ is closed. Thus if π is injective we obtain an inverse $*$ homomorphism $\pi^{-1} : \pi(A) \rightarrow A$, and both must be contractible, hence isometries. \square

Proposition 10.22. *If B is a C^* subalgebra of A , and \mathfrak{a} a closed ideal of A , then $B + \mathfrak{a}$ is a C^* subalgebra of A .*

Proof. Completeness is a three space property, and since \mathfrak{a} is complete, we need only show that $(B + \mathfrak{a})/\mathfrak{a}$ is complete, and from this we will obtain completeness. Let $\pi : A \rightarrow A/\mathfrak{a}$ be the quotient map. But the composed $*$ morphism ψ defined by

$$B \rightarrow A \rightarrow A/\mathfrak{a}$$

tells us that the image of B is closed, and

$$\psi(B) = (B + \mathfrak{a})/\mathfrak{a}$$

so $(B + \mathfrak{a})/\mathfrak{a}$ is complete. \square

10.2.3 Positive Operators

Every C^* algebra is essentially a ring of bounded operators on a Hilbert space. The challenge is to construct a canonical Hilbert space from a C^* algebra. The study of a certain subclass of operators will be essential. A *positive* operator $\phi : A \rightarrow \mathbb{C}$ on a C^* algebra maps A_+ into $[0, \infty)$. Certainly the set of positive maps form a real subspace of A^* .

Example. *If $\phi \in \Phi_A$, then ϕ is positive, for*

$$\phi(M^*M) = |\phi(M)|^2$$

*and every positive element can be written in the form M^*M .*

Example. If X is locally compact and Hausdorff, and μ is a complex Borel measure on X , then the functional

$$f \mapsto \int f d\mu$$

from $C_0(X)$ to \mathbf{C} is positive if and only if μ is a positive measure.

Example. Consider the trace map $\text{tr} : M_n(\mathbf{C}) \rightarrow \mathbf{C}$. Any positive matrix M can be, by a change of basis, put into the form

$$\begin{pmatrix} \lambda_1 & 0 & \dots & 0 \\ 0 & \lambda_2 & \dots & 0 \\ 0 & 0 & \ddots & 0 \\ 0 & 0 & \dots & \lambda_n \end{pmatrix}$$

where $\lambda_1, \dots, \lambda_n \geq 0$. The trace is invariant of a change in basis, and thus

$$\text{tr}(M) = \sum \lambda_i \geq 0$$

Thus the trace is positive.

Example. For a fixed $x \in H$, the map $T \mapsto \langle Tx, x \rangle$ is positive.

One need not verify continuity for positive operators, for we obtain this automatically.

Proposition 10.23. A positive operator $\phi : A \rightarrow \mathbf{C}$ is continuous.

Proof. We claim the supremum

$$K = \sup\{\phi(M) : M \in A_+, \|M\| \leq 1\}$$

is finite. Otherwise, we may pick a sequence M_1, M_2, \dots with $\|M_i\| \leq 1$ and $\phi(M_i) \geq 4^n$. Let

$$M = \sum \frac{M_i}{2^i}$$

Then for each n , $2^n M \geq M_n$, so

$$\phi(M) \geq \frac{\phi(M_n)}{2^n} \geq \frac{4^n}{2^n} = 2^n$$

which yields an immediate contradiction.

Now suppose that M is an arbitrary operator, choose self adjoint N and L in A such that

$$M = N + iL = N_+ - N_- + iL_+ - iL_-$$

Then if $\|M\| \leq 1$,

$$|\phi(M)| = |\phi(N_+) - \phi(N_-) + i\phi(L_+) - i\phi(L_-)| \leq 4K$$

which shows that ϕ is continuous. □

Theorem 10.24. *The following are equivalent*

1. ϕ is positive operator.
2. For every BAI E_α contained within the positive identity on A ,

$$\|\phi\| = \lim_{\alpha} \phi(E_\alpha)$$

3. There is a BAI E_α in the positive identity on A for which

$$\|\phi\| = \lim_{\alpha} \phi(E_\alpha)$$

Proof. Without loss of generality, assume $\|\phi\| = 1$. Suppose ϕ is positive. Fix an approximate identity E_α . Then

$$\limsup_{\alpha} \phi(E_\alpha) \leq 1$$

On the other hand, when $\|M\| \leq 1$,

$$|\phi(M)|^2 = \lim_{\alpha} |\phi(E_\alpha M)|^2 \leq \liminf_{\alpha} \phi(E_\alpha^2) \phi(M^* M)$$

□

Chapter 11

Spectral Theory for Normal Operators

11.1 Resolutions of the Identity

If X is a measure space and H is a Hilbert space, then a *resolution of the identity* on X is a mapping E from measurable subsets of X to self-adjoint projections $B(H)$, such that $E(\emptyset) = 0$, $E(X) = I$, such that $E(S_1 \cap S_2) = E(S_1)E(S_2)$, if S_1 and S_2 are disjoint, $E(S_1 \cup S_2) = E(S_1) + E(S_2)$, and for any $x, y \in H$, the map $E_{x,y}(S) = \langle E(S)x, y \rangle$ is a complex measure. If X is a topological space equipped with the measure space structure given by Borel sets, then we will implicitly assume that each $E_{x,y}$ is a regular Borel measure.

A consequence of these assumptions is that all projections in the image of E commute with one another, and that if S_1 and S_2 are disjoint, $E(S_1)$ and $E(S_2)$ are orthogonal to one another. E is certainly finitely additive, but except in trivial situations, is *not* countably additive in the topology of $B(H)$, since projections have norm one. On the other hand, the measure *is* countably additive in the strong operator topology on $B(H)$, i.e. for each $x \in H$, if $\{S_n\}$ are disjoint, and $S = \bigcup S_n$, then

$$\sum_{n=1}^{\infty} E(S_n)x = E(S)x,$$

because the series converges absolutely, i.e. by the orthogonality of the

projections $E(S_n)$,

$$\sum_{n=1}^{\infty} \|E(S_n)x\|^2 \leq \|x\|^2.$$

We will use resolutions of the identity to obtain a very general functional calculus on operators.

If E is a resolution to the identity, we let $f : X \rightarrow \mathbf{C}$ be a measurable function. The *essential range* of f is the complement of the set formed from the union of all open discs $D \subset \mathbf{C}$ with $E(f^{-1}(D)) = 0$. The function f is essentially bounded if its range is bounded, and we can then define $\|f\|_{L^\infty(X,E)}$ to be the maximum modulus of an element of the essential range of f . The quotient of the family of essentially bounded functions by the closed ideal of functions f with $\|f\|_{L^\infty(X,E)} = 0$ is a Banach algebra denote by $L^\infty(X,E)$. In this Banach algebra, $\sigma(f)$ is precisely the essential range of f . We now define a way to integrate elements of $L^\infty(X,E)$ against the resolution E . The integral will not only be linear, but even *multiplicative*. Recall that a closed normal subalgebra of $B(H)$ is a commutative subalgebra closed under taking adjoints.

Theorem 11.1. *For any resolution of the identity E , there exists an isometric map $L^\infty(X,E) \rightarrow B(H)$ onto a closed normal subalgebra A of $B(H)$, such that if we denote the image of $f \in L^\infty(X,E)$ under this map by*

$$\int_X f dE,$$

then for any $x, y \in H$,

$$\left\langle \left(\int_X f dE \right) x, y \right\rangle = \int f dE_{x,y},$$

and

$$\left\| \left(\int_X f dE \right) x \right\|^2 = \int_X |f|^2 dE_{x,x}.$$

An operator $M \in B(H)$ commutes with every $E(S)$ if and only if M commutes with every element of A , and

$$\|f\|_{L^\infty(X,E)} = \sup_{\|x\| \leq 1} \left(\int_X |f|^2 dE_{x,x} \right)^{1/2}.$$

Proof. TODO: Not that interesting. \square

The spectral theorem for normal operators on a Hilbert space says that each normal $T \in B(H)$ induces a canonical resolution E on Borel subsets of $\sigma(T)$ which allows us to reconstruct T . We can do things in a more general light, by simultaneously performing the spectral theorem for any subalgebra of $B(H)$, all of whose elements are normal.

Theorem 11.2. *Suppose that A is a closed normal unital subalgebra of $B(H)$, and let Φ_A be the maximal ideal space of A . Then there exists a unique resolution E on the Borel sets of Φ_A such that*

$$T = \int_{\Phi_A} \hat{T} dE,$$

where \hat{T} is the Gelfand transform of T . The inverse of the Gelfand transform from $C(\Phi_A)$ to A extends to an isometric $*$ isomorphism Φ of $L^\infty(X, E)$ onto a closed subalgebra B of $B(H)$ containing A , where the inverse is given by

$$f \mapsto \int_{\Phi_A} f dE.$$

The subalgebra B is the closure (in the topology of $B(H)$) of all finite linear combinations of the projections in the image of E . Moreover, if U is a nonempty open subset of Φ_A , then $E(U) \neq 0$. An operator $S \in B(H)$ commutes with every element of A if and only if S commutes with every projection in the image of E .

Proof. TODO \square

In particular, if $T \in B(H)$ is a normal operator, then there exists a unique resolution of the identity on Borel subsets of $\sigma(T)$ such that

$$T = \int_{\sigma(T)} \lambda dE(\lambda)$$

and if S commutes with T , every projection in the image of E commutes with S . We can extend E to a resolution on \mathbb{C} by setting $E(S) = 0$ if $S \cap \sigma(T) = \emptyset$. We call E in this case the *spectral decomposition* of T . To obtain this result from the more general result, we consider the smallest closed subalgebra A of $B(H)$ containing T , T^* and the identity. Then A is a normal (in fact, a commutative) algebra, and Φ_A can be identified

with $\sigma(T)$ in such a way that $\hat{T}(\lambda) = \lambda$ for each $\lambda \in \sigma(T)$. In this case, the inverse map from $L^\infty(\sigma(T), E)$ to $B(H)$ is denoted by $f \mapsto f(T)$. In this special case, we find that $\|f(T)\| \leq \|f\|_{L^\infty(X)}$ (defined by a pointwise supremum rather than an essential supremum), where this is an equality if $f \in C(\sigma(T))$, since then $\|f\|_{L^\infty(X, E)} = \|f\|_{L^\infty(X)}$ by virtue of the fact that if $U \subset \mathbf{C}$ is open, then $E(U) \neq 0$. Intuitively, we can think of $dE(\lambda)$ as giving an ‘infinitesimal projection’ onto the eigenspace corresponding to the eigenvalue λ . Let us consider some applications of this symbolic calculus.

Theorem 11.3. *If $T \in B(H)$ is normal, then*

$$\|T\| = \sup_{\|x\| \leq 1} |\langle Tx, x \rangle|.$$

Proof. If scalars in $\sigma(T)$ always corresponded to eigenvectors, then this result would follow from the fact that $\|T\| = \max \{|\lambda| : \lambda \in \sigma(T)\}$. But elements of $\sigma(T)$ do not always correspond to eigenvectors. On the other hand, we will use the spectral theorem to produce ‘almost eigenvectors’ for T , which will suffice for our purposes. Since T is normal, we know that $\|T\|$ is equal to the spectral radius of T . Thus for any $\lambda_0 \in \sigma(T)$ with $|\lambda_0| = \|T\|$, let B_ε be a ball of radius ε centred at λ_0 . Then if E is the resolution associated with T , we know that $E(B_\varepsilon) \neq 0$. Since $E(B_\varepsilon)$ is a projection, we may find $x \in H$ such that $E(B_\varepsilon)x = x$ and such that $\|x\| = 1$. If $f(\lambda) = (\lambda - \lambda_0)\mathbf{I}_{B_\varepsilon}$, then

$$f(T) = (T - \lambda_0)E(B_\varepsilon)$$

and so $f(T)x = Tx - \lambda_0x$. Now $\|f\|_{L^\infty} \leq \varepsilon$, so it follows that $\|f(T)x\| \leq \varepsilon$, which implies that

$$|\langle Tx, x \rangle - \lambda_0| = |\langle f(T)x, x \rangle - \lambda_0| \leq \|f\|_{L^\infty} \leq \varepsilon.$$

This implies the required result. \square

Remark. If $T : \mathbf{C}^2 \rightarrow \mathbf{C}^2$ is given by $T(a, b) = (b, 0)$, then $\langle T(a, b), (a, b) \rangle = ab$, and so $|\langle Tx, x \rangle| \leq 1/2$ if $\|x\| \leq 1$. But $\|T\| = 1$.

Similar arguments, left as an exercise, show that a normal operator $T \in B(H)$ is self-adjoint if and only if $\sigma(T) \subset \mathbf{R}$, and unitary if and only if $\sigma(T)$ is contained on the unit circle. This need not be true for general operators (e.g. non orthogonal projections are not self-adjoint).

Theorem 11.4. *Let $T \in B(H)$ be normal, and let E be its spectral decomposition. If $f \in C(\sigma(T))$ and $S = f^{-1}(0)$, then the kernel of $f(T)$ is equal to the range of $E(S)$.*

Proof. If $g = \mathbf{I}_S$, then $f \cdot g = 0$, and so $f(T)g(T) = f(T)E(S) = 0$, which shows the range of $E(S)$ is contained in the kernel of $f(T)$. Conversely, suppose $x \in H$ and $f(T)x = 0$. For each $\varepsilon > 0$, let

$$g_\varepsilon(x) = \frac{1}{f(x)} \mathbf{I}_{S_\varepsilon^c}(x).$$

Since f is continuous, $g_\varepsilon \in L^\infty(\sigma(T))$ for all $\varepsilon > 0$. And $f \cdot g_\varepsilon = \mathbf{I}_{S_\varepsilon^c}$, so that $E(S_\varepsilon^c)x = g_\varepsilon(T)f(T)x = 0$. Countable additivity shows that $E(S^c)x = 0$, which shows by orthogonality of $E(S)$ and $E(S^c)$ that x is in the range of $E(S)$. \square

It follows that if $\lambda \in \sigma(T)$, and $E = E(\lambda)$, then the range of E is precisely the kernel of $T - \lambda_0$, i.e. the set of eigenvectors for T with eigenvalue λ . Since $E(U) \neq 0$ for any open U , any isolated point of $\sigma(T)$ is an eigenvalue for some eigenvector. If $\sigma(T)$ is countable, $\sigma(T) = \{\lambda_1, \lambda_2, \dots\}$ then we have, in the strong operator topology, an orthogonal decomposition $I = \sum_{n=1}^\infty E_i$, where $E_i = E(\lambda_i)$, and $TE_i = \lambda E_i$.

Previously, we saw that for a compact normal operator T , $\sigma(T)$ was countable with no limit point but the origin, and $T - \lambda$ had a finite dimensional kernel for each $\lambda \in \sigma(T)$. Using the spectral theorem, we can now prove that this characterizes the compact normal operators from the family of all normal operators.

Theorem 11.5. *A normal operator $T \in B(H)$ is compact if and only if $\sigma(T)$ has no limit point but the origin (which implies $\sigma(T)$ is countable), and for any $\lambda \neq 0$, then $T - \lambda$ has finite dimensional kernel.*

Proof. To show this is sufficient, if $f_\varepsilon(x) = x \cdot \mathbf{I}(|x| \geq \varepsilon)$, then $f_\varepsilon(T)$ is a finite sum of finite rank operators. Thus $f_\varepsilon(T)$ is compact for each $\varepsilon > 0$. On the other hand,

$$\|T - f_\varepsilon(T)\| \leq \|x - f_\varepsilon\|_{L^\infty} \leq \varepsilon,$$

which shows T is the limit of a family of compact operators, and is therefore compact. \square

It follows that if T is normal and compact, then T has an eigenvalue λ with $\|T\| = |\lambda|$, and $f(T)$ is compact for any $f \in C(\sigma(T))$ with $f(0) = 0$, but $f(T)$ is not compact if H is infinite dimensional and $f(0) \neq 0$.

Theorem 11.6. *If $T \in B(H)$, then $\langle Tx, x \rangle \geq 0$ for each $x \in H$ if and only if T is self-adjoint, and $\sigma(T) \subset [0, \infty)$. We call such an operator positive semidefinite and write $T \geq 0$. If $\sigma(T) \subset (0, \infty)$, we say T is positive definite, and then $\langle Tx, x \rangle > 0$ for $x \neq 0$, and if λ_0 is the smallest value in $\sigma(T)$, then $\langle Tx, x \rangle \geq \lambda_0 \|x\|^2$ (and λ_0 is the largest constant for which this bound holds), in which case we write $T \geq \lambda_0$.*

Proof. Suppose $\langle Tx, x \rangle \geq 0$ for each $x \in H$. Then because $\langle Tx, x \rangle$ is real-valued, $\langle Tx, x \rangle = \langle x, Tx \rangle = \langle T^*x, x \rangle$, so that $\langle (T - T^*)x, x \rangle = 0$ for all $x \in H$. Since $T - T^*$ is normal, it follows that $T - T^* = 0$, from which it follows that T is self adjoint. For any $\lambda < 0$, for $x \neq 0$ we have

$$-\lambda \|x\|^2 = \langle -\lambda x, x \rangle \leq \langle (T - \lambda)x, x \rangle \leq \|(T - \lambda)x\| \|x\|,$$

from which it follows that $\|(T - \lambda)x\| \geq -\lambda \|x\|$, so $T - \lambda$ is an isomorphism, hence $\lambda \notin \sigma(T)$.

Conversely, if T is self adjoint, and $\sigma(T) \subset [0, \infty)$, then

$$\langle Tx, x \rangle = \int_0^\infty \lambda dE_{x,x}.$$

Since $E_{x,x}$ is a positive measure, and $\lambda \geq 0$, it follows that $\langle Tx, x \rangle \geq 0$. If $\sigma(T) \subset (0, \infty)$, and λ_0 is the smallest element of $\sigma(T)$, then $\sigma(T - \lambda_0) \subset [0, \infty)$, and is self adjoint, so $\langle (T - \lambda_0)x, x \rangle \geq 0$, which gives $\langle Tx, x \rangle \geq \lambda_0 \|x\|^2$. \square

If $T \in B(H)$ is positive-semidefinite, then we can let $f(t) = \sqrt{t}$ be a well defined continuous function on $\sigma(T)$, and so $f(T) = \sqrt{T}$ is well defined. Thus every positive semidefinite function has a square root. Now $\sigma(f(T)) = \sqrt{\sigma(T)} \subset [0, \infty)$, and since f is real valued, \sqrt{T} is self adjoint. Thus \sqrt{T} is a positive-semidefinite square root for T . This is actually the unique such positive semidefinite square root, because if $S \geq 0$ is a positive semidefinite square root for T , then we let A be the smallest closed subalgebra of $B(H)$ containing S , T , and \sqrt{T} . Then A is commutative, and is isomorphic to $C(\Phi_A)$ for some compact Φ_A , in such a way that positive-semidefinite elements of A are mapped onto non-negative functions in

$C(\Phi_A)$. But since non-negative functions have unique square roots, it follows that $S = \sqrt{T}$.

For any operator T , T^*T is a positive semidefinite self-adjoint operator, so we can consider the operator $|T| = \sqrt{T^*T}$. Then

$$\| |T|x \|^2 = \langle \sqrt{T^*T}x, \sqrt{T^*T}x \rangle = \langle T^*Tx, x \rangle = \|Tx\|^2.$$

This is the *only* positive semidefinite operator with this property, since if $\|Sx\| = \|Tx\|$ for each $x \in H$, then $\langle (S^2 - T^*T)x, x \rangle = 0$ for all x , from which it follows that $S^2 = T^*T$.

Over the complex numbers, any $z \in \mathbf{C}$ can be written as rw for some $r \geq 0$ and $|w| = 1$. This leads to the problem of trying to factorize an operator T into a product UP , where $P \geq 0$ and U is unitary. We call UP a *polar decomposition* of T .

Chapter 12

Unbounded Operators

We often wish to describe a spectral theory of operators that are not necessarily bounded. For instance, in the spectral analysis of partial differential equations with non constant coefficients, one would like to consider a Banach space of functions X supporting both derivative operators $D : X \rightarrow X$ given formally by $Df = f'$, and multiplier operators $M : X \rightarrow X$ given formally by $Mf = xf$. Applying the product rule would show that $[D, M] = DM - MD$ is the identity operator. But the following intriguing result shows that we *cannot* find a Banach space of functions X upon which this theory works.

Lemma 12.1. *There does not exist a unital normed algebra A supporting $x, y \in A$ with $xy - yx = 1$.*

Proof. If $x, y \in A$ were found with $xy - yx = 1$, then it would follow that $x^n y - yx^n = nx^{n-1}$ (easily proved e.g. by induction). But this means that

$$n\|x^{n-1}\| = \|x^n y - yx^n\| \leq 2\|x^n\|\|y\| \leq 2\|x^{n-1}\|\|x\|\|y\|,$$

so that $\|x\|\|y\| \geq n/2$ for every integer n , which is impossible. □

Thus we cannot find a common function space upon which to analyze differentiation and multiplication as bounded operators. Fortunately, we can obtain a theory of *unbounded operators* which allow us to view D and M as operators with a common domain, e.g. $L^2(\mathbf{R})$. The key to this theory is to allow operators that are only defined on dense subsets of Banach spaces.

Let X and Y be Banach spaces. An *unbounded operator* $T : X \rightarrow Y$ is a linear map $T : \mathcal{D}(T) \rightarrow Y$ defined on a subspace $\mathcal{D}(T) \subset X$, which we call the *domain* of the operator. Two unbounded operators are equal if they have the same domain, and are equal as functions. We do not assume this operator has any topological properties in general - if the operator was bounded the Hahn-Banach theorem would imply that it was just the restriction of a bounded operator $T : X \rightarrow Y$, so there would be no real importance to consider the operator on a smaller subspace.

Example. Consider the operator Δ on $L^2(\Omega)$, where Ω is some domain in \mathbf{R}^d which, for simplicity, we might as well assume has smooth boundary. One domain we could take for Δ is $C^\infty(\Omega)$, or the space of all smooth, periodic functions and these will be different unbounded operators. Of more relevance in what follows, we can consider Δ either as an operator with domain \mathcal{D}_D (the space of functions with Dirichlet boundary conditions) consisting of all functions $f \in C^2(\Omega)$ vanishing on the boundary, or the domain \mathcal{D}_N of functions $f \in C^2(\Omega)$ such that $\nabla f \cdot \nu$ vanishes on the boundary, where ν is the normal vector to the boundary. Thus we obtain two different unbounded operators, Δ_D , and Δ_N . The spectral theory of operators on the two domains can certainly differ, e.g. for $\Omega = [0, 1]$, the eigenfunctions of Δ_D are precisely scalar multipliers of the functions

$$f_n(x) = \sin(2\pi nx)$$

which have eigenvalues $-4\pi^2 n^2$, for $n > 0$. The eigenfunctions for Δ_N , on the other hand, are scalar multiples of the functions

$$f_n(x) = \cos(2\pi nx)$$

which also have eigenvalues $-4\pi^2 n^2$, for $n \geq 0$. Thus 0 is an eigenvalue for Δ_N , but not for Δ_D . The differences between Neumann and Dirichlet boundary conditions only get more distinct in higher dimensions.

For any unbounded operator $T : X \rightarrow Y$, we can consider the graph $\Gamma(T) = \{(x, Tx) : x \in \mathcal{D}(T)\}$. If $\Gamma(T)$ is a closed subset of $X \times Y$, we say that T is *closed*. The closed graph theorem implies that if T is closed and $\mathcal{D}(T) = X$, then T is actually a bounded operator. For any closed operator T , we can give $\mathcal{D}(T)$ the structure of a Banach space by defining

$$\|x\|_T = \sqrt{\|x\|^2 + \|Tx\|^2},$$

and in this topology T will then be a bounded operator from $\mathcal{D}(T)$ to Y .

The importance of being closed is *essential* to the spectral theory of unbounded operators. Given an unbounded operator $T : X \rightarrow Y$, we define the resolvent $\rho(T)$ of T to be the set of all complex numbers $\lambda \in \mathbb{C}$ such that $T - \lambda$ has a bounded everywhere defined inverse $S : X \rightarrow \mathcal{D}(T)$, i.e. such that $(T - \lambda)S = I_X$, and $S(T - \lambda) = I_{\mathcal{D}(T)}$, which we will denote by $(T - \lambda)^{-1}$. The spectrum $\sigma(T)$ is then the complement of the resolvent of T .

Lemma 12.2. *If $\sigma(T) \neq \mathbb{C}$, then T is closed.*

Proof. Suppose $\lambda \notin \sigma(T)$, and assume without loss of generality that $\lambda = 0$. Let $S : X \rightarrow \mathcal{D}(T)$ be the inverse of T . If $x_i \rightarrow x$ in $\mathcal{D}(X)$, and $Tx_i \rightarrow y$ in X , then by continuity, $x_i = STx_i \rightarrow Sy$, so $x = Sy$. But this means that $Tx = TSy = y$. Thus T is closed. \square

Like with bounded operators, $\sigma(T)$ is a closed set, but it need not be compact. The proof is roughly the same.

Lemma 12.3. *If $\lambda \in \rho(T)$, and $L^{-1} = \|(T - \lambda)^{-1}\|$, then $B(\lambda, L) \subset \rho(T)$.*

Proof. Assume without loss of generality that $\lambda = 0$. Let $S = T^{-1}$. Then

$$\begin{aligned} (T - \lambda)^{-1} &= \sum_{n=0}^{\infty} (-\lambda)^n \\ (T - \lambda)^{-1} &= \sum_{n=0}^{\infty} (-\lambda)^n T^n \end{aligned}$$

\square

We wish to extend the natural operations on the class of operators to the class of unbounded operators. Given a scalar λ , and an unbounded operator $T : X \rightarrow Y$, we can define an unbounded operator $\lambda T : X \rightarrow Y$ with $\mathcal{D}(\lambda T) = \mathcal{D}(T)$. For two unbounded operators $T : \mathcal{D}(T) \rightarrow Y$ and $S : \mathcal{D}(S) \rightarrow Y$, we can define $T + S$ as an unbounded operator defined on $\mathcal{D}(T) \cap \mathcal{D}(S)$. Similarly, we can define $T - S$. If $T : \mathcal{D}(T) \rightarrow Y$ and $S : \mathcal{D}(S) \rightarrow Z$ are unbounded operators, we can define ST as an unbounded operator defined on $\mathcal{D}(T) \cap T^{-1}(\mathcal{D}(S))$. The distributive law $(S + U)T = ST + UT$ continues to hold since

$$\mathcal{D}(T) \cap T^{-1}(\mathcal{D}(S) \cap \mathcal{D}(U)) = [\mathcal{D}(T) \cap T^{-1}(\mathcal{D}(S))] \cap [\mathcal{D}(T) \cap T^{-1}(\mathcal{D}(U))],$$

but the distributive law $T(S+U) = TS+TU$ might *not* hold in this scheme. If $T : X \rightarrow Y$ is an injective unbounded operator, we can define $T^{-1} : Y \rightarrow X$ to be the unbounded operator whose domain is the image of T .

An unbounded operator is *densely defined* if $\mathcal{D}(T)$ is a dense subset of X . Given an unbounded densely defined operator $T : \mathcal{D}(X) \rightarrow Y$, we can define an unbounded operator T^* with domain

$$\mathcal{D}(T^*) = \{\gamma \in Y^* : \gamma \circ T \in \mathcal{D}(T)^*\},$$

where $\mathcal{D}(T)$ is given the subspace topology inherited from X . We then set $T^*\gamma$ to be the unique $\lambda \in X^*$ extending the continuous functional $\gamma \circ T$. In inner product notation, we then have

$$\langle T^*\gamma, x \rangle = \langle \gamma, Tx \rangle$$

for any $x \in \mathcal{D}(T)$ and $\gamma \in \mathcal{D}(T^*)$. Note that if T is not densely defined, this extension might not be unique, and so the adjoint is not well defined. If $T : X \rightarrow Y$ and $S : Y \rightarrow Z$ are densely defined unbounded operators such that ST is also densely defined, then $(ST)^*$ is an extension of T^*S^* . If S is bounded, then the domains are actually equal to one another, and so we get the familiar identity $T^*S^* = (ST)^*$. We also note that if $T : X \rightarrow Y$ is densely defined and unbounded, and $S : Y^* \rightarrow X^*$ is another unbounded operator such that $\gamma(Tx) = (S\gamma)(x)$ for any $x \in \mathcal{D}(T)$ and $\gamma \in \mathcal{D}(S)$, then $\mathcal{D}(S) \subset \mathcal{D}(T^*)$, because $\gamma \circ T$ is a bounded functional on $\mathcal{D}(T)$ for any $\gamma \in \mathcal{D}(S)$.

If $T : H_1 \rightarrow H_2$ is an unbounded operator between two Hilbert spaces H_1 and H_2 , the unbounded adjoint can be naturally identified with an unbounded adjoint $T^* : H_2 \rightarrow H_1$, which we will do whenever we are talking about Hilbert spaces. If $T : H \rightarrow H$ is unbounded, and T^* is an extension of T , then we say T is *symmetric*. This means that for any $x_1, x_2 \in \mathcal{D}(T)$, $\langle Tx_1, x_2 \rangle = \langle x_1, Tx_2 \rangle$. If we actually have $T^* = T$, which means that in addition to the previous property, $x \mapsto \langle Tx, y \rangle$ is continuous if and only if $y \in \mathcal{D}(X)$, then we say T is *self adjoint*. However, it is possible for T to be a strict extension of T^* , as the next example shows, and we do not say this operator is symmetric nor self-adjoint.

Example. Let $T : L^2[0, 1] \rightarrow L^2[0, 1]$ be the unbounded operator whose domain consists of all absolutely continuous functions $f : [0, 1] \rightarrow \mathbb{C}$ such that $f' \in L^2[0, 1]$, and for such functions, let

$$Df = if'.$$

If $g \in L^2[0,1]$ satisfies

$$\left| \int f'(x)g(x) dx \right| \lesssim \|f\|_{L^2[0,1]}$$

for all $f \in \mathcal{D}(T)$, then we claim that g is absolutely continuous, and $g(0) = g(1) = 0$. Conversely, if g is absolutely continuous, then

$$\int f'(x)g(x) dx = f(1)g(1) - f(0)g(0) - \int f(x)g'(x) dx,$$

and we see that this quantity is $\lesssim \|f\|_{L^2[0,1]}$ precisely when $g(0) = g(1) = 0$. All that remains is to argue that g is absolutely continuous, but this follows because $\mathcal{D}(T)$ contains the restriction of $C_c^\infty(0,1)$, so the boundedness equation above implies that the distributional derivative g' of g in $\mathcal{D}(0,1)^*$ is square integrable, which implies that g lies in $\mathcal{D}(T)$, i.e. g is absolutely continuous and $g' \in L^2[0,1]$. Thus we see that, even though we have $\langle Tx, y \rangle = \langle x, Ty \rangle$ for $x \in \mathcal{D}(T)$ and $y \in \mathcal{D}(T^*)$, T^* is only defined on a proper subset of the domain of T . On the other hand, the operator $S = T^*$, which extends T , is symmetric, but not self adjoint, since one can verify that $S^* = T^{**} = T$, and that S^* therefore extends S . The operator T' defined as the restriction of T to functions f with $f(0) = f(1)$, is self adjoint, since one can easily verify from the integration by parts formula that $\mathcal{D}((T')^*) = \mathcal{D}(T')$.

We now come to an interesting result which relates the graph of T and T^* .

Theorem 12.4. *Let T be a densely defined unbounded operator on a Hilbert space H . Then $\Gamma(T^*) = [\Gamma(T)]^\perp$, where $V(x, y) = (-y, x)$, and $H \times H$ is viewed as a Hilbert space. In particular, T^* is always a closed unbounded operator, all self-adjoint operators are closed.*

Proof. The following four statements are equivalent to one another:

- $(y, z) \in \Gamma(T^*)$
- $\langle y, Tx \rangle = \langle z, x \rangle$ for every $x \in \mathcal{D}(T)$.
- $\langle (y, z), V(x, Tx) \rangle = 0$ for every $x \in \mathcal{D}(T)$.
- (y, z) lies in $V(\Gamma(T))^\perp$.

Thus the proof is complete. \square

Here are some corollaries for a symmetric unbounded operator $T : H \rightarrow H$:

- If $\mathcal{D}(T) = H$, then T is a bounded, self-adjoint operator.
- If T is injective and self-adjoint, then it has dense image, and T^{-1} is self adjoint.
- If the image of T is dense, then T is injective.
- If T is surjective (so that (c) applies and T is actually injective), then T is self-adjoint, and T^{-1} is bounded.

Proof. To see (a), if $\mathcal{D}(T) = H$, then because T^* extends T , $T^* = T$, and since T^* is closed, it follows that T is bounded and self-adjoint.

To obtain (b), to show the image is dense let x be perpendicular to the image of T . Then $x \in \mathcal{D}(T^*)$ because $\langle x, Tx' \rangle = 0$ for all $x' \in \mathcal{D}(T)$, and is thus trivially continuous. But this implies that $x \in \mathcal{D}(T^*) = \mathcal{D}(T)$, and $Tx = T^*x = 0$. Thus $x = 0$, which shows the image of T is dense. Thus T^{-1} is a densely defined unbounded operator, and so we can consider its adjoint. We claim that $[V\Gamma(T^{-1})]^\perp = \Gamma(T^{-1})$, which would show that T^{-1} is self adjoint. But $V(\Gamma T^{-1}) = \Gamma(-T)$ and $V\Gamma(-T)$, and since T is self-adjoint, $-T$ is self adjoint, and so $\Gamma(-T)^\perp = V\Gamma(-T) = \Gamma(T^{-1})$, which completes the argument.

To prove (c), we suppose that $Tx = 0$. Then for any $x' \in \mathcal{D}(T)$, $\langle x, Tx' \rangle = \langle Tx, x' \rangle = 0$. Thus x is orthogonal to the image of T . But this implies $x = 0$.

To prove (d), we note that T^{-1} is symmetric, because if $x_1, x_2 \in H$ then $x_i = Ty_i$ for some $y_i \in \mathcal{D}(T)$, and so applying the fact that T is symmetric, we find that

$$\langle T^{-1}x_1, x_2 \rangle = \langle y_1, Ty_2 \rangle = \langle Ty_1, y_2 \rangle = \langle x_1, T^{-1}x_2 \rangle.$$

Since T^{-1} is self adjoint, injective, and has dense image, this implies $T = (T^{-1})^{-1}$ is self adjoint. But (a) also implies that T^{-1} is bounded. \square

A powerful consequence of this result is that if T is a closed operator, then we have an orthogonal decomposition

$$H \times H = V\Gamma(T) \oplus \Gamma(T^*)$$

Thus for any $a, b \in H$, there exists a unique pair $(x, y) \in \mathcal{D}(T) \times \mathcal{D}(T^*)$ such that $a = y - Tx$ and $b = x + T^*y$. Moreover, for this pair we have

$$\|a\|^2 + \|b\|^2 = \|x\|^2 + \|Tx\|^2 + \|y\|^2 + \|T^*y\|^2.$$

In particular, setting $a = 0$, we find that with any closed, unbounded operator T , we can associate a pair of bounded operators $A_1 : H \rightarrow \mathcal{D}(T)$ and $A_2 : H \rightarrow \mathcal{D}(T^*)$, with $\|A_1\| \leq 1$ and $\|A_2\| \leq 1$, such that $TA_1 = A_2$, $I = A_1 + T^*A_2$, and for any $x \in H$,

$$\|x\|^2 = \|A_1x\|^2 + \|TA_1x\|^2 + \|A_2x\|^2 + \|T^*A_2x\|^2.$$

In particular, if $Q = I + T^*T$, then $I = QA_1$, so Q is surjective onto H . Moreover, Q is also injective, since if $x \in \mathcal{D}(Q)$, $Tx \in \mathcal{D}(T^*)$, and so

$$\langle x, x \rangle + \langle Tx, Tx \rangle = \langle x, x \rangle + \langle T^*Tx, x \rangle = \langle Qx, x \rangle$$

which shows by Cauchy-Schwartz that $\|x\|^2 \leq \|x\|\|Qx\|$, so that $\|Qx\| \geq \|x\|$, and thus Q is injective. Thus Q is a bijection from $\mathcal{D}(Q)$ to H , which means that I is an extension of A_1Q . The fact that $\langle Qx, x \rangle \geq 0$ for all $x \in \mathcal{D}(Q)$, means $\langle A_1x, x \rangle \geq 0$ for all $x \in H$, i.e. A_1 is positive semidefinite. Since A_1 is injective and self-adjoint, this implies $Q = A_1^{-1}$ is self-adjoint. And this means $T^*T = Q - I$ is self adjoint for any closed unbounded operator T .

Theorem 12.5. *If $T : H \rightarrow H$ is a densely defined closed unbounded operator, then $\mathcal{D}(T^*)$ is dense, and $T^{**} = T$.*

Proof. We use the fact that $H \times H = V\Gamma(T) \oplus \Gamma(T^*)$. Suppose $y \perp \mathcal{D}(T^*)$. Then $(y, 0)$ is perpendicular to $\Gamma(T^*)$. This means that $(y, 0) \in V\Gamma(T)$, so $(0, -y) \in \Gamma(T)$, which implies $y = 0$. Thus $\mathcal{D}(T^*)$ is dense. To show $T^{**} = T$, we simply note that we can also write $H \times H = V\Gamma(T^{**}) \oplus \Gamma(T^*)$, from which it follows that $V\Gamma(T^{**}) = V\Gamma(T)$, and so $T = T^{**}$. \square

A symmetric operator T is *maximally symmetric* if it has no proper symmetric extension. Any self-adjoint operator T is maximally symmetric, for if S is symmetric and extends T , then $T = T^*$ extends S^* , which means S extends S^* , and since S is symmetric, S^* extends S , hence we conclude $S = S^*$, and as a result, $S = T$. On the other hand, maximally symmetric operators need not be self-adjoint.

12.1 The Cayley Transform

One trick to understanding symmetric operators T is to consider the operator $T + i$. This is because if T is symmetric, a simple calculation shows that for $x \in \mathcal{D}(T)$,

$$\|Tx + ix\|^2 = \|x\|^2 + \|Tx\|^2.$$

Thus $\Gamma(T)$ is isometric to the image of $T + i$. Thus $T + i$ is injective, T is closed if and only if the range of $T + i$ is closed, and if $T + i$ is onto H , then T is maximally symmetric, because otherwise $T + i$ would have a symmetric extension which was still injective, which is impossible. Exactly the same results are true if we replace i with any other purely imaginary quantity.

The mapping

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

is a one-to-one-correspondence between \mathbf{R} and the unit circle in \mathbf{C} (except for one). It follows from the continuous functional calculus that the map

$$T \mapsto (T - i)(T + i)^{-1}$$

is a one-to-one-correspondence between the family of self-adjoint bounded operator T on a Hilbert space H , and the family of unitary operators on H which do not contain 1 in their spectrum. We wish to work with ‘unbounded’ isometries, which will extend the theory of the Cayley transform to unbounded symmetric operators. An unbounded isometry is an unbounded operator $U : H \rightarrow H$ such that $\|Ux\| = \|x\|$ for each $x \in \mathcal{D}(U)$. Equivalently, this holds if and only if $\langle Ux, Uy \rangle = \langle x, y \rangle$ for each $x, y \in \mathcal{D}(U)$. If U is an unbounded isometry, and $I - U$ has dense range, then $I - U$ is injective, for if $Ux = x$, then for any $y \in \mathcal{D}(U)$,

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

Because the range of $I - U$ is dense, $x = 0$. It is similarly easy to prove that the three sets $\mathcal{D}(U)$, $U(\mathcal{D}(U))$, and $\Gamma(U)$ are all closed, or all not closed.

Let us now show the Cayley transform gives a correspondence between symmetric unbounded operators on H , and unbounded isometries on H (unbounded operators $U : H \rightarrow H$ such that $\|Ux\| = \|x\|$ for each $x \in \mathcal{D}(U)$), or equivalently, $\langle Ux, Uy \rangle = \langle x, y \rangle$ for $x, y \in \mathcal{D}(U)$). For any symmetric operator T , since $\|Tx + ix\|^2 = \|x\|^2 + \|Tx\|^2$, the map $U = (T - i)(T + i)^{-1}$ is an unbounded isometry with domain is the image of the map $T + i$. The following properties then holds

Theorem 12.6. Suppose U is the Cayley transform of T .

- U is closed if and only if T is closed.
- The map $I - U$ is injective from $\mathcal{D}(U)$ onto $\mathcal{D}(T)$, and $T = i(I + U)(I - U)^{-1}$.
- U is unitary if and only if T is self-adjoint.

Conversely, any unbounded isometry on H such that $I - U$ is injective is the Cayley transform of a symmetric operator.

Proof. TODO

□

If U_1 and U_2 are Cayley transforms of symmetric operators T_1 and T_2 , then T_1 extends T_2 if and only if U_1 extends U_2 . Thus problems about extensions of symmetric operators reduce to problems about extensions of isometries. Given an unbounded symmetric closed operator T , the range of $T + i$ and $T - i$ are closed, and U is an isometry from the image of the first map to the image of the second map. The dimensions of $(T + i)^\perp$ and $(T - i)^\perp$ are called the *deficiency indices* of the map T . Then it follows from the properties above that T is self-adjoint if and only if both of its deficiency indices are zero, T is maximally symmetric if and only if at least one of its deficiency indices is zero, and T has a self-adjoint extension if and only if its two deficiency indices are equal.

Example. Let V be the right shift operator on l^2 . Then V is an isometry, and $I - V$ is injective. Thus V is the Cayley transform of a symmetric operator T on l^2 . Since V has domain l^2 , and its range has codimension one, the deficiency indices of T are zero and one. Thus T is an example of a maximally symmetric operator which is not self adjoint.

12.2 The Spectral Theory of Unbounded Symmetric Operators

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12.3 Operator Semigroups

TODO

Chapter 13

Von Neumann Algebras

13.1 Spectral Theorem for Normal Operators

In finite dimensional theory, a normal linear transformation $T : \mathbb{C}^n \rightarrow \mathbb{C}^n$ can be written

$$T = \sum \lambda_i P_i$$

where P_i is projection onto the eigenspace corresponding to the eigenvalue λ_i . In this section we apply the theory of C^* algebras to extend this theorem to arbitrary normal bounded transformations from a Hilbert space to itself. The trick is to ‘integrate’ rather than sum up the subspaces.

A *spectral measure* on a measure space (Ω, \mathcal{F}) is a function $E : \mathcal{F} \rightarrow B(H)$, for which

$$E(\Omega) = \text{id}_H$$

such that $E(S)$ is a projection for each S , and such that for a disjoint family of sets $\{S_1, S_2, \dots\}$ in \mathcal{F} , then

$$E\left(\bigcup S_i\right) = \sum E(S_i)$$

with convergence pointwise (the strong operator topology), rather than in the operator norm. It follows that $E(\emptyset) = 0$.

Lemma 13.1. *If $S \subset W$, $E(S) \leq E(W)$.*

Proof. Since $E(W - S)$ is a projection, $E(W - S) \geq 0$, so

$$E(W) = E(W - S) + E(S) \geq E(S)$$

Hence a spectral measure is monotone on projections. □

Lemma 13.2. *If S and W are disjoint sets, then $E(S)E(W) = 0$.*

Proof. $E(S \cup W) = E(S) + E(W)$, so

$$E(S \cup W) = E(S \cup W)E(S \cup W) = E(S) + E(W)E(S) + E(S)E(W) + E(W)$$

which implies

$$E(S) + E(W) = E(S) + E(W)E(S) + E(S)E(W) + E(W)$$

hence $E(W)E(S) + E(S)E(W) = 0$. □

Lemma 13.3. $E(S \cap W) = E(S) \circ E(W)$.

Proof. We have

$$E(S \cap W) + E(S - W) = E(S) \quad E(S \cap W) + E(W - S) = E(W)$$

and

$$E(S \cup W) = E(S \cap W) + E(S - W) + E(W - S)$$

From which it follows that

$$\begin{aligned} E(S) + E(W) &= 2E(S \cap W) + E(S - W) + E(W - S) \\ &= E(S \cap W) + E(S \cup W) \end{aligned}$$

Now multiply both sides of the equation on the right by $E(W)$ gives

$$E(S)E(W) + E(W)^2 = E(S \cap W)E(W) + E(S \cup W)E(W)$$

Since $S \subset S \cup W$, $E(S) \leq E(S \cup W)$, and

$$E(S \cup W)E(S) = E(S)$$

$$\text{Now } E(S \cup W)E(S) = E(S)^2 + E(W - S)E(S) = E(S) \quad \square$$

Example. *If $H = \mathbf{C}$, then $B(H) \cong \mathbf{C}$, and any $\{0, 1\}$ valued measure μ is a spectral measure, provided $\mu(\Omega) = 1$. If $\mu(S_1) = \mu(S_2) = 1$, then*

$$1 = \mu(S_1) = \mu(S_1 \cap S_2) + \mu(S_1 - S_2)$$

$$1 = \mu(S_2) = \mu(S_1 \cap S_2) + \mu(S_2 - S_1)$$

If $\mu(S_1 \cap S_2) = 0$, then $\mu(S_1 - S_2) = \mu(S_2 - S_1) = 1$, which implies

$$\mu((S_1 - S_2) \cup (S_2 - S_1)) = \mu(S_1 - S_2) + \mu(S_2 - S_1) = 2$$

an impossibility. Thus $\mu(S_1 \cap S_2) = 1$. If $\mu(S_1) = 0$, then $S_1 \cap S_2 \subset S_1$, so $\mu(S_1 \cap S_2) \leq \mu(S_1) = 0$. The same holds if $\mu(S_2) = 0$. Thus we have verified that $\mu(S_1 \cap S_2) = \mu(S_1) \circ \mu(S_2)$ on a case by case basis. The countable summation property holds by the property of the measure itself.

Example. If H is infinite dimensional, and $N \in K(H)$, then

$$\sigma(N) = \{\lambda_0, \lambda_1, \dots\}$$

where $\lambda_0 = 0$, and each λ_i is an eigenvalue. For each $n > 0$, let P_n be orthogonal projection onto $\ker(\lambda_n - N)$, and let P_0 project onto the orthogonal complement of the closed linear space of the images of P_1, P_2, \dots . For $S \subset \sigma(N)$, let

$$E(S) = \sum_{\lambda_n \in S} P_n$$

In the sense that the sum on the right is interpreted pointwise. Then E is a spectral measure on $(\sigma(N), \mathcal{P}(\sigma(N)))$. Surely

$$E(\sigma(N)) = \sum_{n=0}^{\infty} P_n = id_H$$

And since $P_i \circ P_j = 0$ if $i \neq j$,

$$\begin{aligned} E(S_1) \circ E(S_2) &= \left(\sum_{\lambda_n \in S_1} P_n \right) \circ \left(\sum_{\lambda_m \in S_2} P_m \right) \\ &= \sum_{\lambda_n \in S_1, \lambda_m \in S_2} P_n P_m \\ &= \sum_{\lambda_n \in S_1 \cap S_2} P_n = E(S_1 \cap S_2) \end{aligned}$$

And the countable additivity follows by the pointwise definition of the sum of operators.

Example. If Ω is a σ -finite measure space with measure μ , and $H = L^2(\Omega)$, define

$$E(S)\xi = \chi_S \xi$$

Then E is a spectral measure. Surely $E(\Omega) = id_H$, $E(\emptyset) = 0$, because $\chi_\Omega = 1$, $\chi_\emptyset = 0$. Now $E(S_1 \cap S_2) = E(S_1) \circ E(S_2)$ follows because $\chi_{S_1 \cap S_2} = \chi_{S_1} \chi_{S_2}$. Similarly, $\chi_{\cup S_i} = \sum \chi_{S_i}$ pointwise if the S_i are disjoint.