

My Report about Operator Theory

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Chapter 1

Introduction

The most of this report is based on the two books, A Course in Functional Analysis and A Course in Operator Theory by John B. Conway. However, I reorganized the contents under my understanding and added some other materials to help me understand the structure of the knowledge of the fundamental operator theory.

My understanding of this structure is based on "General-to-special" order. Like, we firstly research the general space, and then by adding more other structure on it, we get a more special space with more interesting properties. Conversely, these properties also revealed the essential traits of the more special space. Then we can strengthen the traits the special space has, thus a more special space is born. Using this order, the appearance of C^* -algebra and von Neumann algebra is very nature for me.

Firstly, I have seen the structure of finite dimensional linear space, but how to research the infinite dimensional linear space is a question. Clearly, the algebraic structure cannot provide enough information about the infinite dimension. However, by learning topological spaces, I know that the topology can provide a method to research the infinite sets, thus equipping the infinite dimensional linear space with an appropriate topology is a good ideal. Naturally, this topology should be compatible with the original linear structure. Therefore, topological vector spaces is important, and there are many interesting properties generated by the merge of the topology and the algebra. If the topology of the T.V.S. can be more special, like the pre-"length" of elements is continuous with respect to the topology, then we get a locally convex topological vector space with more interesting properties than the general T.V.S. has. When learning the finite dimensional linear space, I saw the corresponding dual space is very "like" with the original space. But how do these properties change in the infinite dimensional case? The Hahn-Banach Theorem provides an answer.

Now, the locally convex T.V.S. can be more special. If there is only one seminorm generating the locally convex T.V.S., how can the space become?

In order to make the space be Hausdorff, the seminorm need to a norm. Moreover, since completed topological space is more interesting, this special space should be completion. Therefore, we get a Banach space. And then all of above properties the general T.V.S. and locally convex T.V.S. have can be applied to the Banach space and have more powerful form. Moreover, the space consisting of all continuous maps between the Banach spaces can be equipped with more interesting structure, which can also be a Banach space with the induced norm.

The norm generating a Banach space can be more special. If the norm, which satisfies the Parallelogram Law, can be generated by an inner product, then the Banach space with this inner product is a Hilbert space. Because of this extra inner product structure, the Hilbert space has the Projection Theorem, and thus the dual space is as same as the original space, and thus the the space consisting of all continuous maps between the same Hilbert space can contain the adjoint version of each of them. This space is the operator space.

For researching the operator space, similarly, we can firstly research a more general version. The main difference between the operator space and the Banach space is that the operator space can do multiplication, defined as composition. Moreover, this multiplication is continuous with respect to the norm topology. Therefore, the general version of the operator space is a Banach space with a extra continuous multiplication, then it becomes a Banach algebra. Comparing the Banach space, the Banach algebra has one more algebraic structure, and the invertibility is an important property with respect to this structure. Thus, the appearance of the concept of spectrum is naturally. Besides, the commutative Banach algebra can have an explicit expression.

Besides, the operator algebra has one more operation than the general Banach algebra have, the involution. And the involution is not only continuous with respect to the norm topology. In fact, it provides an identity. That means the topology is determined by the algebraic structure. Then, the general version is C^* -algebra, whose topology is completely determined the algebraic structure it has. Because of the strong relation between the topology and the algebra, C^* -algebra has some similar properties to the complex number system. By the GNS construction, any C^* -algebra can be embedded in a operator algebra. It is amazing, since the general version of operator algebra is itself. Thus, researching operator algebras is enough.

However, the C^* -subalgebras of an operator algebra may be too small to contain enough "interesting" elements. Thus, finding a weaker topology and extending the C^* -algebra to a larger subalgebra are good ideals. But how can we do is? The Continuous Functional Calculus can provide a clue. We have known that the continuous function space on a compact is weakly dense in the bounded Borel function space and the Gelfand transform maps the continuous function space to the operator algebra. So it is natural to

find a method of extending the map from the continuous function space to the bounded Borel function space. In general, we should define the operator-valued measure, thus the integral of operator makes sense. Comparing this progress to the progress embedding the continuous function space in the bounded Borel function space, the new topologies we need appeared naturally. That is the weak operator topology and the strong operator topology. Then the new algebra is born, von Neumann algebra. And in fact, these "interesting" elements are projections. And since the norm topology is determined by the algebraic structure, it can also use the algebraic properties to construct von Neumann algebras. Similarly, the structure of an abelian von Neumann algebra may be easier. There is a particular abelian von Neumann algebra, which is generated by the normal operator. By extending the Gelfand transform, the structure of this von Neumann algebra can become concrete. Moreover, any abelian von Neumann algebra has same structure.

From finite rank operators to infinite rank operators, there is a special class of operators, which can be approximated by a sequence of finite rank operators, compact operators. Because of it, compact operators has a lot of properties as similar as finite rank operators, such as the spectrum, the representations and so on. There are also some interesting classes contained in the compact operators, like the trace class. The importance of trace class is that it is the dual space of the operator algebra, and thus it can provide an extra topology on the operator algebra. This topology has some relations with the weak operator topology, and then it can provide a new method to research the von Neumann algebra. By the quotient map, there is a weak version compact operators, the Fredholm operators. So the Fredholm theory appeared. Finally, we have known that any normal operator can be diagonalizable in the finite dimensional case, but in the infinite dimensional case, we need the compact perturbation in order to get similar result.

For the general von Neumann algebra, combining the *WOT* and the ultraweak topology, we get more properties than a general C^* -algebra has, like the structure of the ideals. Moreover, it provides the equivalent condition the simple von Neumann algebra has, which is called factor. We have known that the von Neumann algebra is generated by all of its projections, thus these projections play an important role. By researching the relations of projections and the properties they have, it can classify von Neumann algebras into corresponding classes. In fact, these projections can be classified as abelian ones, non-abelian ones and pure infinite ones, and the corresponding von Neumann algebras are Type I, Type II and III. Type I is easier, but Type II and II do not have trivial example.

Chapter 2

Topological Vector Spaces and Banach Spaces

2.1 Topological Vector Spaces

Linear operations, i.e addition and scalar multiplication, provide an algebraic structure on a set, therefore constructing a vector space. In the course, linear algebra, we have learn the algebraic structure of finite dimensional vector spaces. But how to deal with infinite dimensional vector spaces? By learning the topological spaces, we know that the topological structure can give us a method to research properties of infinity. Thus, we need to equip a vector space with an additional topological structure, which should coincide with the algebraic structure. This is the reason why we define the topological vector space.

2.1.1 Topological Spaces

■ **Definition:** First, we define the topological structure on a general set.

Definition 1. A topological space $X = (X, \mathcal{T}_X)$ consists of a set X , called the underlying space of X , and a family \mathcal{T}_X of subsets of X s.t.

- 1) $X, \emptyset \in \mathcal{T}_X$.
- 2) if $U_\alpha \in \mathcal{T}_X$ for $\alpha \in A$, then $\bigcup_{\alpha \in A} U_\alpha \in \mathcal{T}_X$.
- 3) if $U_1, U_2 \in \mathcal{T}_X$, then $U_1 \cap U_2 \in \mathcal{T}_X$.

And, \mathcal{T}_X is called a topology on X . The element in \mathcal{T}_X is called open set.

Thus, the topological structure on a set X is totally determined by the family \mathcal{T}_X . In particular, from 2), we can simplify \mathcal{T}_X . In other words, like the basis of a vector space, there is a "basis" of \mathcal{T}_X .

Definition 2. If X is a set, a basis for a topology on X is a family \mathcal{B} of subsets of X , s.t.

- 1) $\forall x \in X, \exists B \in \mathcal{B}, \text{ s.t. } x \in B.$
- 2) if $x \in B_1 \cap B_2$, where $B_1, B_2 \in \mathcal{B}$, then there is a $B_3 \in \mathcal{B}$, s.t. $B_3 \subset B_1 \cap B_2.$

\mathcal{B} can generate a topology \mathcal{T}_X on X by doing infinite times union of elements in \mathcal{B} . In fact,

$$\mathcal{T}_X = \{ U \subset X : U = \bigcup_{\alpha \in A} U_\alpha, U_\alpha \in \mathcal{B} \}, \text{ where } A \text{ is any index set.}$$

We can do more on the basis by 3).

Definition 3. A subbasis of a topology on X is a collection \mathcal{S} of subsets of X , whose union is X . The topology generated by \mathcal{S} is noted by $\mathcal{T}(\mathcal{S})$.

In fact, we can generate the basis of $\mathcal{T}(\mathcal{S})$ by \mathcal{S} , that is,

$$\mathcal{B} = \{ U \subset X : U = \bigcap_{\alpha=1}^n U_\alpha, U_\alpha \in \mathcal{S}, n \in \mathbb{N} \}$$

Therefore, this basis can generate the coincided topology on X .

■ **Continuous maps:** Next, we need to endow the general maps with the topological structure.

Definition 4. Let X, Y be topological spaces and $f: X \rightarrow Y$ be a map. We say f is continuous if

$$\forall V \in \mathcal{T}_Y, f^{-1}(V) \in \mathcal{T}_X.$$

In other words, if $f: X \rightarrow Y$ is a continuous map, \mathcal{T}_X has "more" elements than \mathcal{T}_Y has. To get more rigorous discription, we define the following concept.

Definition 5. Let X be a set and $\mathcal{T}, \mathcal{T}'$ be two topologies on X we say that \mathcal{T} is coarser than \mathcal{T}' if $\mathcal{T} \subset \mathcal{T}'$.

Therefore, if $f: X \rightarrow Y$ is a continuous map, the topology $\mathcal{T}(f^{-1}(\mathcal{T}_Y))$ is coarser than \mathcal{T}_X .

■ **Generating topologies:** We use above methods to generate some interested topologies on a set.

- 1) Initial topology: Given maps $f_\alpha: X \rightarrow Y_\alpha$ ($\alpha \in A$) from a set X to a family of topological spaces $\{Y_\alpha: \alpha \in A\}$. Let

$$\mathcal{S} = \{ f_\alpha^{-1}(V): V \in \mathcal{T}_{Y_\alpha}, \alpha \in A \}$$

Then $\mathcal{T}(\mathcal{S})$ is the coarsest topology on X such that each f_α is continuous, called the initial topology induced by the family of maps $\{f_\alpha: \alpha \in A\}$.

- 2) Final topology: Given maps $f_\alpha: X_\alpha \rightarrow Y$ ($\alpha \in A$) from a family of topological spaces $\{X_\alpha: \alpha \in A\}$ to a set Y . Let

$$\mathcal{S} = \{ V: f_\alpha^{-1}(V) \in \mathcal{T}_{X_\alpha}, \alpha \in A \}$$

Then $\mathcal{T}(\mathcal{S})$ is the finest topology on Y such that each f_α is continuous, called the final topology induced by the family of maps $\{f_\alpha: \alpha \in A\}$.

Here is some important examples using above way to generate topologies.

Example 6. (Initial topology)

- 1) Subspace topology: Let Y be a topological space and $X \subset Y$ be a subset of Y . The inclusion map $i: X \rightarrow Y$ can generate the initial topology \mathcal{T}_X on X , then X become a subspace of Y . In fact,

$$\mathcal{T}_X = \{ X \cap U: U \in \mathcal{T}_Y \}$$

- 2) Product topology: Let $\{Y_\alpha: \alpha \in A\}$ be a family of topological spaces. The product set is

$$\prod_{\alpha \in A} Y_\alpha = \{ A \xrightarrow{f} \bigcup_{\alpha \in A} Y_\alpha: \forall \alpha \in A, f(\alpha) \in Y_\alpha \}$$

There is a family of maps $\{p_\beta: \prod_{\alpha \in A} Y_\alpha \rightarrow Y_\beta (\beta \in A)\}$. Therefore, these maps can generate the initial topology \mathcal{T} on $\prod_{\alpha \in A} Y_\alpha$ and let $\prod_{\alpha \in A} Y_\alpha$ be the product topological space. In fact,

$$\mathcal{T} = \{ \prod_{\alpha \in A} V_\alpha: V_\alpha \in \mathcal{T}_{Y_\alpha} \text{ \& } \#\{\alpha \in A: V_\alpha \neq Y_\alpha\} < \infty \}$$

Remark. The condition, $\#\{\alpha \in A: V_\alpha \neq Y_\alpha\} < \infty$, is because that when the subbasis generates the basis, only finite many elements can do intersection.

Example 7. (Final topology)

Quotient topology: For a topological space X on which an equivalent relation R is fixed, $\pi: X \rightarrow X/R$ is the quotient map, then the quotient set can be equipped with the final topology \mathcal{T} generated by the quotient map. Therefore, X/R become a quotient topological space. In fact,

$$\mathcal{T} = \{ U \subset X/R: \pi^{-1}(U) \in \mathcal{T}_X \}$$

▪ **Countability and metrizable topology:** When learning the analysis of real functions, we usually use the sequence to describe the topological properties. But in some general topology, we cannot just use sequence since some properties of "uncountability". In this case, the concept of net can be applied to some "uncountable" topologies. Furthermore, there is a class of more special topology, metrizable topology, which has some better properties.

Definition 8. (Net)

1) Direct set: A direct set (D, \geq) consists of a nonempty set D and a relation \geq on D , satisfies:

- i) $\forall d \in D, d \geq d$
- ii) $\forall d_1, d_2, d_3 \in D, \text{ if } d_3 \geq d_2 \ \& \ d_2 \geq d_1, \text{ then } d_3 \geq d_1$
- iii) $\forall d, d' \in D, \exists d'' \text{ s.t. } d'' \geq d \ \& \ d'' \geq d'.$

2) if X is a set, a net is a map $x: D \rightarrow X$ from a direct set D to X

Example 9. If X is a topological space and $x \in X$, then let

$$D = \{ \text{all open neighbourhoods of } x \}, \ U \geq V \Leftrightarrow U \subset V$$

Then D is a direct set and $x_\alpha (\alpha \in D)$ is a net. And we say $x_\alpha \rightarrow x$ if and only if

\forall open neighbourhood U of x in $X, \exists \delta \in D, \forall \alpha \in D$ with $\alpha \geq \delta \Rightarrow x_\alpha \in U$

Nets can be used as sequences in topological spaces. Like,

Proposition 10. If X is a topological space and the net $x_\alpha (\alpha \in D)$ defined above and $A \subset X$, then

- 1) $\overline{A} = \{ x \in X : \exists x_\alpha \text{ in } A, x_\alpha \rightarrow x \}$
- 2) $f: X \rightarrow Y$ is continuous between two topological spaces, $x_0 \in X, f$ is continuous at x_0 , if and only if
 \forall net $x_\alpha (\alpha \in D), \text{ s.t. } x_\alpha \rightarrow x_0 \Rightarrow f(x_\alpha) \rightarrow f(x_0)$

Definition 11. (Countability)

- 1) First countability: For a topological space X, X is called first countable if for each point $x \in X, x$ has a countable neighbourhood basis.
- 2) Second countability: A topological space X is second countable if it has the countable topological basis.

Remark. Clearly, the second countable topological space is first countable, but the converse is not true.

In particular, if X is first countable, sequences can be used to illuminate topological properties rather than nets. Like,

Proposition 12. *If X is first countable, then*

- 1) $U \subset X$ is closed $\Leftrightarrow \forall x \in U, \exists$ a sequence $\{x_n\} \subset U$, s.t. $x_n \rightarrow x$.
- 2) sequential compactness is equivalent to compactness.

And for the second countability, it is about the separability.

Definition 13. (Separability)

- 1) A subset A of a topological space X is called dense if $\overline{A} = X$.
- 2) A topological space is called separable if it has a countable dense subset.

By the definition, we can clearly know that:

Proposition 14. *If X is a second countable topological space, then it is separable and every open covering of X has a finite subcollection covering X .*

We can classify topological spaces into some classes.

Definition 15. X is a topological space, then we call X is:

- (T_0) $\forall x, y \in X, \exists$ open $U \subset X$, s.t. $x \in U$ but $y \notin U$ or $y \in U$ but $x \notin U$
(Kolmogorov space)
- (T_1) $\forall x, y \in X, \exists$ open $U, V \subset X$, s.t. $x \in U$ but $y \notin U$ and $y \in V$ but $x \notin V$
($\Leftrightarrow \forall x \in X, \{x\}$ is closed)
- (T_2) $\forall x, y \in X, \exists$ open $U, V \subset X$, s.t. $x \in U$ & $y \in V$ and $U \cap V = \emptyset$
(Hausdorff space)
- (T_3) T_1 holds and $\forall x \in X$ and closed $C \subset X$, if $x \notin C$, then \exists open $U, V \subset X$, s.t. $x \in U$ & $C \subset V$ and $U \cap V = \emptyset$ (regular space)
- (T_4) T_1 holds and \forall closed $C_1, C_2 \subset X$, if $C_1 \cap C_2 = \emptyset$, then \exists open $U, V \subset X$, s.t. $C_1 \subset U$ & $C_2 \subset V$ and $U \cap V = \emptyset$
(normal space)

Then we can specify a class of more powerful topological space.

Definition 16. If X is a topological space, then X is said to be metrizable if there exists a metric d on the set X that induces the topology of X .

Remark. Clearly, if X is metrizable, X is second countable and normal.

Here is two metrization theorems provides the essence of metric spaces.

Theorem 17. (*Metrization theorems*)

Urysohn A topological space is separable and metrizable if and only if it is regular, Hausdorff and second countable.

Nagata–Smirnov A topological space is metrizable if and only if regular, Hausdorff and has a σ -locally finite basis.

■ **Complete metric space:** For a metric space, we know it is first countable, so the concept of net is unnecessary. And thus sequences are enough to determine the topological structures, like that sequential compactness is equivalent to compactness.

Proposition 18. A compact subset of a metric space is closed, bounded and separable.

Remark. it is clearly, since compactness is also about finity.

For any metric space, we can use the following theorem to get a completion of that and this completion is unique. Thus, we can always assume a metric space is complete.

Theorem 19. Let (X, d) be a metric space. Then, there exists a metric space (\hat{X}, \hat{d}) with the following properties:

- 1) (\hat{X}, \hat{d}) is complete.
- 2) There is an embedding σ from X to \hat{X} .
- 3) $\sigma(X)$ is dense in \hat{X} .

And this (\hat{X}, \hat{d}) is unique with respect to isomorphism.

Complete metric space is important since it is "sufficiently large". Rigorously, we can use the following definition to describe it.

Definition 20. (Baire Category) A metric space is said to be of the first category if it can be written as a countable union of sets that are nowhere dense. Otherwise, it is of the second category.

Proposition 21. A complete metric space is a space of the second category.

■ **Filters:** For convenience, we define some terminologies.

Definition 22. A filter on a set X is a family \mathcal{F} of subsets of X satisfying the following conditions:

- 1) $\emptyset \notin \mathcal{F}$
- 2) \mathcal{F} is closed under finite many intersections

- 3) Any subset of X containing a set in \mathcal{F} belongs to \mathcal{F} .

Example 23. For a topological space X and $x \in X$, and let

$$\mathcal{F}(x) = \{ \text{all neighbourhoods of } x \}$$

Then $\mathcal{F}(x)$ is a filter and $\mathcal{F}(x)$ satisfies the following properties:

- 1) $\forall U \in \mathcal{F}(x), x \in U$
- 2) $\forall U \in \mathcal{F}(x), \exists V \in \mathcal{F}(x), \text{ s.t. } \forall y \in V, U \in \mathcal{F}(y)$

And conversely, if we can find $\mathcal{F}(x)$ for any $x \in X$ with above two properties, these can define a unique topology \mathcal{T} s.t. $\mathcal{F}(x)$ is the filter of neighbourhoods of x for any $x \in X$. In fact,

$$\mathcal{T} = \{ U \subset X : x \in U \Rightarrow U \in \mathcal{F}(x) \}$$

Also, we can define the basis of $\mathcal{F}(x)$, noted by $\mathcal{B}(x)$. That is $\mathcal{B}(x) \subset \mathcal{F}(x)$ with the following properties:

- 1) $\forall U \in \mathcal{B}(x), x \in U$
- 2) $\forall U_1 \& U_2 \in \mathcal{B}(x), \exists U_3 \in \mathcal{B}(x), \text{ s.t. } U_3 \subset U_1 \cap U_2$
- 3) If $y \in U \in \mathcal{B}(x), \exists W \in \mathcal{B}(y), W \subset U$

2.1.2 Definition and Properties

■ **Definition:** Now, we need to endow the topological structure on a vector spaces. And the most important thing is that the topological structure should coincide with the algebraic structure.

Definition 1. A vector space X over a field \mathbb{K} (where $\mathbb{K} = \mathbb{C} \text{ or } \mathbb{R}$) is called a topological vector space if X is equipped with a topology \mathcal{T} s.t. the addition and the scalar multiplication, i.e.

$$\begin{aligned} (x, y) &\mapsto x + y \\ (\lambda, x) &\mapsto \lambda x \end{aligned}$$

are continuous with respect to the topology \mathcal{T} .

In this definition, the most important part is that the addition and the scalar multiplication are continuous. This condition provides some additional properties for the topology and also for the linear operations. First, it can simply the topology.

Proposition 2. Given a t.v.s. X ,

- 1) For any $x_0 \in X$, the map $x \mapsto x + x_0$ is a homeomorphism.

2) For any $\lambda \in \mathbb{K}$, then map $x \mapsto \lambda x$ is a homeomorphism.

Proof. It is clearly, since by the definition, $x \mapsto x - x_0$ and $x \mapsto \frac{1}{\lambda}x$ are continuous. \square

Therefore, the topology of a t.v.s is completely determined by the filter of neighbourhoods of any point. Or, more rigorously,

Corollary 3. For a t.v.s X , the filter $\mathcal{F}(x)$ of neighbourhoods of $x \in X$ is as same as $\{ U + x : U \in \mathcal{F}(e) \}$, where e is the unit element in X .

Thus, to research the topology of a t.v.s. X , we just need to research the filter $\mathcal{F}(e)$ of neighbourhoods of e . First, there are two special properties of some subsets of a t.v.s. X .

Definition 4. For a subset U of a t.v.s. X ,

- 1) U is absorbing if $\forall x \in X, \exists \rho > 0$ s.t. $\forall \lambda \in \mathbb{K}$ with $|\lambda| \leq \rho$, we have $\lambda x \in U$.
- 2) U is balanced if $\forall x \in U, \forall \lambda \in \mathbb{K}$ with $|\lambda| \leq 1$, we have $\lambda x \in U$.

Then, the following theorem reveals the essence of $\mathcal{F}(e)$.

Theorem 5. A filter \mathcal{F} of a vector space X over \mathbb{K} is the filter of neighbourhoods of the unit element e w.r.t. some topology compatible with the algebraic structure of X if and only if

- 1) $\forall U \in \mathcal{F}, e \in U$
- 2) $\forall U \in \mathcal{F}, \exists V \in \mathcal{F}$ s.t. $V + V \subset U$
- 3) $\forall U \in \mathcal{F}, \forall \lambda \in \mathbb{K}$ with $\lambda \neq 0, \lambda U \in \mathcal{F}$
- 4) $\forall U \in \mathcal{F}, U$ is absorbing
- 5) $\forall U \in \mathcal{F}, \exists V \in \mathcal{F}$ s.t. $V \subset U$ is balanced

Proof. If $\mathcal{F} = \mathcal{F}(e)$, these statements clearly hold.

1) is trivial.

2) is true since the addition is continuous.

3) and 4) hold since the scalar multiplication is continuous.

For 5), because the scalar multiplication is continuous, we can find a $W \in \mathcal{F}$ s.t. $\lambda W \subset U$ for any $|\lambda| \leq \rho$, then let $V = \bigcup_{|\lambda| \leq \rho} \lambda W$. Clearly, $V \in \mathcal{F}$ and V is balanced.

Conversely, We can define

$$\mathcal{F}(x) = \{ U + x : U \in \mathcal{F} \}$$

for any $x \in X$. It can be easily checked that $\mathcal{F}(x)$ satisfies the conditions in Example 23 in last subsection. Therefore, these $\mathcal{F}(x)$ can determine a unique topology \mathcal{T} on X .

Now, we just need to check the continuity of the addition and the scalar multiplication. The addition is continuous, since \mathcal{F} satisfies 2). Using conditions 2) and 4) and 5) to get a balanced absorbing open neighbourhood in \mathcal{F} , and this neighbourhood prove the continuity of the scalar multiplication. \square

Here is some simple properties of a t.v.s. X . These properties are directly obtained by definition and above theorem.

Proposition 6. *For a t.v.s. X ,*

- 1) *proper subspaces of X are never absorbing. In particular, if $M \subset X$ is a open subspace, then $M = X$.*
- 2) *each linear subspace of X , endowed with subspace topology, is also a t.v.s.*
- 3) *if H is a linear subspace of X , then \overline{H} is also a linear subspace of X .*
- 4) *if Y is also a t.v.s. and $f: X \rightarrow Y$ is a linear map, then f is continuous if and only if f is continuous at the unit element e .*

■ **Hausdorff t.v.s. :** The Hausdorff Space is important since it can let the concept of limit make sense. And the topology of a t.v.s. can be simplified and has some additional properties, we can get a easier condition that make a t.v.s. become Hausdorff.

Proposition 7. *A t.v.s X is a Hausdorff space if and only if for any $x \in X$ with $x \neq e$ there exists a $U \in \mathcal{F}(e)$ s.t. $x \notin U$.*

Proof. Since the open neighbourhoods of any point in X is completely determined by the open neighbourhoods of e , this proposition is equivalent to the statement that (T_1) implies Hausdorff.

The proof can be accomplished by obtaining a contradiction to the given condition that $x \neq e$, $\exists U \in \mathcal{F}(e)$ s.t. $x \notin U$. For that U , there is a balanced $V \in \mathcal{F}(e)$ s.t. $V + V \subset U$ and the balance implies that $V - V \subset U$. Therefore, $(x + V) \cap V = \emptyset$. If not, $x + v_1 = v_2$ for $v_1, v_2 \in V$. This implies that $x = v_1 - v_2 \in V - V \subset U$. Thus it is a contradiction. \square

The following theorem is more explicit.

Theorem 8. *For t.v.s. X the following statements are equivalent.*

- 1) *X is Hausdorff.*

2) the intersection of all neighbourhoods of the unit element e is $\{e\}$.

3) $\{e\}$ is closed.

Proof. Before the rigorously proving, the intuition is clearly. Since in a t.v.s. (T_1) is equivalent to Hausdorff, the equivalence of 1) and 3) is clearly true.

- 1) \Rightarrow 2) It is because that elements in $\mathcal{F}(e)$ can separate e and other points.
- 2) \Rightarrow 3) If $x \in \overline{\{e\}}$, i.e. $\forall V_x \in \mathcal{F}(x), V_x \cap \overline{\{e\}} \neq \emptyset \Rightarrow e \in V_x$, and $V_x = U + x$ for some $U \in \mathcal{F}(e)$, then $u + x = e$ for some $u \in U$. Thus, $x = -u \in -U$ for all $U \in \mathcal{F}(e)$. That implies $x = e$.
- 3) \Rightarrow 1) By above mentioned, it just needs to check that if for any topology space Y , $\{y\}$ is closed $\forall y \in Y$, Y is (T_1) .
 Since $\{y_1\}$ is closed, $Y \setminus \{y_1\}$ is open. That means if $y_2 \neq y_1$, there exists a open neighbourhood U of y_2 s.t. $y_1 \notin U$. Similarly, we can find a open neighbourhood V of y_1 s.t. $y_2 \notin V$. Therefore, Y is (T_1) . \square

■ **Quotient t.v.s. :** For a linear subspace M of a t.v.s. X , the quotient topology on X/M can be obtained by the quotient map $\pi: X \rightarrow X/M$. But because of the algebraic structure, it has more properties.

Proposition 9. For a linear subspace M of a t.v.s. X , the quotient map $\pi: X \rightarrow X/M$ is open.

Proof. Let $V \subset X$ be open, then we have

$$\pi^{-1}(\pi(V)) = V + M = \bigcup_{m \in M} (V + m)$$

Since V is open, $V + m$ is open. Thus $\pi^{-1}(\pi(V))$ is open. And by the definition of the topology on X/M , $\pi(V)$ is open. \square

Corollary 10. For a linear subspace M of a t.v.s. X , the quotient space X/M endowed with the quotient topology is a t.v.s..

Proof. We have the following commutative graph, where f and g are corresponding addition maps or scalar multiplication maps on X and X/M .

$$\begin{array}{ccc} X \times X & \xrightarrow{f} & X \\ \downarrow \pi \times \pi & & \downarrow \pi \\ X/M \times X/M & \xrightarrow{g} & X/M \end{array}$$

Then for an open set $V \subset X/M$, since f and π are continuous, and π is open, $(\pi \times \pi) \circ f^{-1} \circ \pi^{-1}(V)$ is open. By above commutative graph, we have $g \circ (\pi \times \pi) = \pi \circ f$. Therefore, $g^{-1}(V)$ is open, i.e. g is continuous. \square

Also, we can find the condition that lets the quotient topological vector space be Hausdorff.

Proposition 11. *Let X be a t.v.s..*

- 1) *M be a linear subspace of X . Then X/M is Hausdorff if and only if M is closed.*
- 2) *$X/\{\overline{e}\}$ is Hausdorff.*

Proof. 2) is true because 1). And 1) clearly holds since M is the unit element in X/M and Theorem 8 in this subsection. \square

Remark. By this method, for any t.v.s., we can find a Hausdorff space w.r.t it.

2.1.3 Continuous Linear Maps

The interesting maps between two topological vector spaces not only preserve the algebraic structure, but also the topological structure, thus these are continuous linear maps.

First, for a linear map $f: X \rightarrow Y$ between vector spaces X and Y , we have the commutative graph, where $\tilde{f}(x + \ker f) = f(x)$ is well-defined.

$$\begin{array}{ccc} X & \xrightarrow{f} & \text{Im } f \xrightarrow{i} Y \\ \downarrow \pi & \nearrow \tilde{f} & \\ X/\ker f & & \end{array}$$

Proposition 1. *Let $f: X \rightarrow Y$ be a linear map between two t.v.s.'s X and Y .*

- 1) *If Y is Hausdorff and f is continuous, then $\ker f$ is closed.*
- 2) *By above notation, f is continuous if and only if \tilde{f} is continuous.*

Proof. 1) is because that $\ker f = f^{-1}(\{e\})$ and Y is Hausdorff. For 2), if \tilde{f} is continuous, it is clearly that $f = i \circ \tilde{f} \circ \pi$ is continuous. Conversely, it is because of the universal property of quotient maps. And in this case, let $U \subset \text{Im } f$ be open, then $f^{-1}(U)$ is open and $\tilde{f}^{-1}(U) = \pi(f^{-1}(U))$. Since π is open, $\tilde{f}^{-1}(U)$ is open. Thus, \tilde{f} is continuous. \square

2.1.4 Complete Topological Vector Spaces

We have just defined the completeness on a metric space by using sequence, but in metric spaces, we know the topology is so powerful that sequences can do any thing, but in general topology, or the topology in a t.v.s., we need an equivalent concept to describe the completeness.

Definition 1. (Completeness)

- 1) A filter \mathcal{F} on a subset A of a t.v.s. X is said to be a Cauchy filter if

$$\forall U \in \mathcal{F}(0) \text{ in } X, \exists M \subset A \text{ s.t. } M \in \mathcal{F} \text{ \& } M - M \subset U$$

- 2) A subset A of a t.v.s. X is said to be complete if every Cauchy filter on A converges to a point $x \in A$.

Remark. Said "the filter converges to a point" means that we can define a net on this filter, and this net converge a point. And this definition is also valid without the algebraic structure.

By this definition, and using the factor that Hausdorff spaces let the limit point of a net uniquely exist, we have similar results comparing with the metric spaces.

Proposition 2. *Let X be a t.v.s..*

- 1) *If X is Hausdorff, any complete set is closed.*
- 2) *If X is complete, any closed set is complete.*

We known any metric space can be completion. Similarly, the same result can obtained in any t.v.s..

Theorem 3. *Let X be a Hausdorff t.v.s., then there exists a complete Hausdorff t.v.s. \hat{X} and a map $i: X \rightarrow \hat{X}$ with the following properties.*

- 1) *i is a topological monomorphism.*
- 2) *$\overline{i(X)} = \hat{X}$.*
- 3) *For any complete Hausdorff t.v.s. Y and for every continuous linear map $f: X \rightarrow Y$, there exists a continuous map $\hat{f}: \hat{X} \rightarrow Y$, s.t. the following graph is commutative*

$$\begin{array}{ccc} X & \xrightarrow{f} & Y \\ \downarrow i & \nearrow \hat{f} & \\ \hat{X} & & \end{array}$$

And (\hat{X}, \hat{f}) is unique with respect to the isomorphism

Proof. The proof is similar as the proof of the completion of metric spaces, which constructs the \hat{X} as a set of equivalent classes of Cauchy sequences. In a t.v.s., we just need to replace Cauchy sequences by Cauchy filters (in fact, Cauchy nets). Let

$$\begin{aligned}
\tilde{X} &= \{ \text{all Cauchy filters in } X \} \\
R &: \mathcal{F} R \mathcal{G} \Leftrightarrow \forall U \in \mathcal{F}(e), \exists A \in \mathcal{F} \text{ \& } B \in \mathcal{G} \text{ s.t. } A - B \subset U \\
\hat{X} &= \tilde{X}/R
\end{aligned}$$

We can easily define linear operations and topology, s.t. \hat{X} become a complete t.v.s.. Then we just need to check the statements in above theorem hold. \square

2.1.5 Finite Dimensional Topological Vector Spaces

For a finite dimensional topological vector space, the topology compatible with the algebraic structure has some properties coincided with the "finiteness". First, continuous linear functionals on a t.v.s. have some properties.

Lemma 1. *Let X be a t.v.s. over \mathbb{K} . Fixed $v \in X$, then the $\phi_v: \mathbb{K} \rightarrow X$ by $\xi \mapsto \xi v$ is continuous,*

Proof. It is because that $\phi_v = f \circ \psi_v$ where f is the multiplication map.

$$\begin{array}{ccccc}
\mathbb{K} & \xrightarrow{\psi_v} & \mathbb{K} \times X & \xrightarrow{f} & X \\
\xi & \mapsto & (\xi, v) & \mapsto & \xi v
\end{array}$$

\square

Lemma 2. *For a non-zero linear functional $L: X \rightarrow \mathbb{K}$, where X is a t.v.s. over \mathbb{K} , the following statements are equivalent.*

- 1) L is continuous,
- 2) $\ker L$ is closed,
- 3) $\ker L$ is not dense in X ,
- 4) L is bounded in some neighbourhood of the origin in X .

Proof. The equivalence of 1) and 2) and 4) is clearly.

- 1) \Rightarrow 2) It is because that $\ker L = L^{-1}(\{0\})$.
- 2) \Rightarrow 3) Since L is non-zero, it clearly holds.
- 3) \Rightarrow 4) By the assumption, there exists a balanced set $V \in \mathcal{F}(e)$ and a point $x \notin \overline{\ker L}$ s.t. $(x + V) \cap \ker L = \emptyset$. $L(V)$ is balanced on \mathbb{K} , therefore $L(V)$ is bounded or $L(V) = \mathbb{K}$. But since $(x + V) \cap \ker L = \emptyset$, $L(V) \neq \mathbb{K}$.
- 4) \Rightarrow 1) This implies that L is continuous at e . But since X is a t.v.s., L is continuous at every point. \square

Theorem 3. *Let X be a finite dimensional Hausdorff t.v.s. over \mathbb{K} (endowed with the standard topology), and $\dim X = d$. Then we have:*

- 1) X is topologically isomorphic to \mathbb{K}^d ,
- 2) every linear functional on X is continuous,
- 3) every linear map from X to any t.v.s. Y is continuous

Proof. For 1), we just need to find a homeomorphic isomorphism from \mathbb{K}^d to X , like the following map, where $\{e_i\}_{i=1}^d$ is the basis of X .

$$\begin{array}{ccc} \mathbb{K}^d & \xrightarrow{\phi} & X \\ (\lambda_1, \lambda_2, \dots, \lambda_d) & \longmapsto & \lambda_1 e_1 + \lambda_2 e_2 + \dots + \lambda_d e_d \end{array}$$

ϕ is clearly an algebraic isomorphism. Thus we just need to check ϕ is both continuous and open.

Check: ϕ is continuous.

When $d = 1$, it is continuous by above lemma. For the general case, since d is finite, ϕ is continuous.

Check: 2) holds and ϕ is open.

When $d = 1$, it is clearly 1) and 2) are true. And suppose 1) and 2) hold for $\dim X \leq d-1$, then when $\dim X = d$, let L be a non-zero linear function on X . Then since $X/\ker L \cong \text{Im } L \subset \mathbb{K}$, $\dim \ker L = d-1$. Therefore, $\ker L \cong \mathbb{K}^{d-1} \Rightarrow \ker L$ is complete $\Rightarrow \ker L$ is closed $\Rightarrow L$ is continuous by above lemma. And,

$$\begin{array}{ccc} X & \xrightarrow{\phi^{-1}} & \mathbb{K}^d \\ \lambda_1 e_1 + \lambda_2 e_2 + \dots + \lambda_d e_d & \longmapsto & (\lambda_1, \lambda_2, \dots, \lambda_d) \end{array}$$

is continuous since each

$$\lambda_1 e_1 + \lambda_2 e_2 + \dots + \lambda_d e_d \longmapsto \lambda_i$$

is continuous.

Then for 3), it is clearly since $\dim \text{Im } L < \infty$. □

Corollary 4.

- 1) Every finite dimensional Hausdorff t.v.s. is complete.
- 2) Every finite dimensional subspace of a Hausdorff t.v.s. is closed.
- 3) For a finite dimensional vector space, there is only one topology w.r.t. homeomorphism that can make it be a Hausdorff t.v.s..
- 4) Every bounded subset on a finite dimensional Hausdorff t.v.s. is compact.

Proof. These properties can be easily obtained by regarding the t.v.s. as \mathbb{K}^d endowed with the standard topology. □

Finally, the most important theorem in this subsection is that the converse of 4) in above corollary is also true.

Theorem 5. *A Hausdorff t.v.s. is locally compact if and only if it is finite dimensional.*

Proof. Let X be a locally compact Hausdorff t.v.s. and K be a compact neighbourhood of e in X , i.e.

$$\exists x_1, \dots, x_r \in X \text{ s.t. } K \subset \bigcup_{i=1}^r (x_i + \frac{1}{2}K)$$

Let $M = \text{span}\{x_1, \dots, x_r\}$, and M is closed. Therefore, X/M is a Hausdorff t.v.s.. Let $\pi: X \rightarrow X/M$ be the quotient map.

Since $K \subset M + \frac{1}{2}K$, $\pi(K) \subset \pi(\frac{1}{2}K)$. Thus, by iterating $\pi(2^n K) \subset \pi(K)$.

As K is absorbing, $X = \bigcup_{n=1}^{\infty} 2^n K$,

$$X/M = \pi(X) = \bigcup_{n=1}^{\infty} \pi(2^n K) \subset \pi(K) \subset X/M$$

And since π is continuous, $\pi(K)$ is compact, i.e. X/M is compact.

claim: $\dim X/M = 0$

Suppose $\dim X/M > 0$, then for some $\overline{x_0} \in X$ with $\overline{x_0} \neq \overline{e}$, $\mathbb{R}\overline{x_0} \subset X/M$.

And since X/M is Hausdorff compact and $\mathbb{R}\overline{x_0}$ is closed, $\mathbb{R}\overline{x_0}$ is compact, which is a contradiction. \square

2.2 Locally Convex Topological Vector Spaces

The locally convex topological vector space is a topological vector spaces whose topology is generated by a family of seminorms, thus it can provide more properties.

2.2.1 Definition by Convex Sets

Firstly, the conception of locally convex space can be obtained by convex sets. So, we need to research some elementary traits of convex sets.

Definition 1.

- 1) Let S be any subset of a vector space X over \mathbb{K} . The convex hull of S , $\text{conv}(S)$, is the smallest convex subset containing S . In fact,

$$\text{conv}(S) = \left\{ \sum_{i=1}^n \lambda_i x_i : x_i \in S, \lambda_i \in [0, 1], \sum_{i=1}^n \lambda_i = 1, n \in \mathbb{N} \right\}$$

- 2) A subset S of a vector space X over \mathbb{K} is absolutely convex, if S is convex and balanced
- 3) A subset S of a vector space X over \mathbb{K} is called a barrel if S is closed, absorbing and absolutely convex.

Proposition 2.

- 1) Arbitrary intersections of convex sets are convex sets, and the sum of two convex sets is convex, and linear maps preserve convex.
- 2) The convex hull of a balanced set is balanced.
- 3) The closure and the interior of a convex set in a t.v.s. is convex.
- 4) Every neighbourhood of the origin of a t.v.s. is contained in a neighbourhood of the origin which is a barrel.

Proof. 1) and 2) are easily obtained by the definition.

To prove 3), for any $\lambda \in [0, 1]$, we define the map

$$\begin{aligned} \phi : X \times X &\longrightarrow X \\ (x, y) &\mapsto \lambda x + (1 - \lambda)y \end{aligned}$$

Using the fact that $\phi(S \times S) \subset S$, $S \times S \subset \phi^{-1}(S)\phi^{-1}(\overline{S})$. And since X is a t.v.s., ϕ is continuous. Thus $\phi^{-1}(\overline{S})$ is closed, i.e. $\overline{S} \times \overline{S} \subset \phi^{-1}(\overline{S}) \Rightarrow \phi(\overline{S} \times \overline{S}) \subset \overline{S}$. \overline{S} is convex.

Since $z + u = \lambda x + (1 - \lambda)y + \lambda u + (1 - \lambda)u$ if $z = \lambda x + (1 - \lambda)y$, $x + U$ & $y + U \subset S \Rightarrow z + U \subset S$. That means \mathring{S} is convex.

For 4), $\forall U \in \mathcal{F}(e)$, the set

$$\overline{\text{conv} \left(\bigcup_{|\lambda| \leq 1} \lambda U \right)}$$

is a barrel. □

Corollary 3. Every neighbourhood of the origin in a t.v.s. X is contained in a neighbourhood of the origin which is absolutely convex.

Now, we can get the definition of a locally convex t.v.s..

Definition 4. A t.v.s. X is said to be locally convex if there is a basis of neighbourhoods of the origin in X consisting of convex sets.

By this definition, the structure of the neighbourhoods of the origin in a locally convex t.v.s. can be more explicit by using above proposition.

Proposition 5. Let X be a locally convex t.v.s..

1) X has a basis of neighbourhoods of origin consisting with open absorbing absolutely convex sets.

2) X has a basis of neighbourhoods of origin consisting with barrels.

Theorem 6. *If X is a locally convex t.v.s., then there exists a basis \mathcal{B} of neighbourhoods of origin consisting of absorbing absolutely convex set s.t.*

$$\forall U \in \mathcal{B}, \forall \rho > 0, \exists W \in \mathcal{B} \text{ s.t. } W \subset \rho U$$

Conversely, if \mathcal{B} is a collection of absorbing absolutely convex subsets of a vector space satisfying above condition, it can generate a unique locally convex t.v.s..

2.2.2 Definition by Seminorms

Definition 1. Let X be a vector space over \mathbb{K} . A map $p: X \rightarrow \mathbb{R}^+$ is called a seminorm if it satisfies:

- 1) $p(x + y) \leq p(x) + p(y), \forall x, y \in X,$
- 2) $p(\lambda x) = |\lambda| p(x), \forall x \in X, \forall \lambda \in \mathbb{K}.$

Remark. In fact, $\ker p$ is a linear subspace and if $\ker p = \{0\}$, p is called a norm.

By the intuition, the seminorm could construct the continuity of the addition and multiplication since it satisfies above properties. Now, we can build this rigorously.

Definition 2. Let X be a vector space and $A \subset X$ be a nonempty subset. The Minkowski functional of A is the map

$$\begin{aligned} p_A : X &\longrightarrow \mathbb{R} \\ x &\longmapsto \inf \{ \lambda > 0 : x \in \lambda A \} \end{aligned}$$

Let X be a vector space and p is a seminorm on X , then let $\mathring{U}_p = \{ x \in X : p(x) < 1 \}$, $U_p = \{ x \in X : p(x) \leq 1 \}$. Thus, U may be the basis generating the topology. To see it, the following proposition is helpful.

Proposition 3. *Let $A \subset X$ be a nonempty subset of a vector space, which is absorbing and absolutely convex, then p_A is a seminorm and $\mathring{U}_{p_A} \subset A \subset U_{p_A}$. Conversely, if q is a norm on X then \mathring{U}_q is an absorbing absolutely convex set and $q = p_{\mathring{U}_q}$.*

Proof. Since A is balanced, $\xi A \in \lambda A \Leftrightarrow x \in \frac{\lambda}{|\xi|} A$. Thus,

$$p_A(x) = |\xi| \inf \left\{ \frac{\lambda}{|\xi|} : x \in \frac{\lambda}{|\xi|} A \right\} = |\xi| p_A(x)$$

And $p_A(x) < \infty$ ($\forall x \in X$) since A is absorbing.

Fixed $x, y \in X, \forall \varepsilon > 0, \exists \lambda, \mu > 0$, s.t. $x \in \lambda A$ & $y \in \mu A$ and

$$\lambda \leq p_A(x) + \varepsilon, \mu \leq p_A(y) + \varepsilon,$$

By convexity of A , $\lambda A + \mu A \subset (\lambda + \mu)A$. Thus,

$$p_A(x) = \inf \{ \delta > 0 : x \in \delta A \} \leq \lambda + \mu \leq p_A(x) + p_A(y) + 2\varepsilon$$

$\Rightarrow p_A(x)$ is a seminorm.

$$\begin{aligned} x \in \mathring{U}_{p_A} &\Rightarrow \exists \lambda \in [0, 1] \text{ s.t. } x \in \lambda A \subset A \\ x \in A &\Rightarrow 1 \in \inf \{ \lambda > 0 : x \in \lambda A \} \Rightarrow p_A(x) \leq 1 \Rightarrow x \in U_{p_A} \end{aligned}$$

That means $\mathring{U}_{p_A} \subset A \subset U_{p_A}$.

Finally, the statements for q can be obtained easily by the definition. \square

Now, we can give the definition of a locally convex t.v.s. by seminorms coinciding with the definition by convex sets.

Theorem 4. *Let X be a vector space and $\mathcal{P} = \{p_i\}_{i \in I}$ be a family of seminorms. Then the initial topology \mathcal{T}_P generated by \mathcal{P} makes X be a locally convex t.v.s. In fact, the basis of neighbourhoods of the origin in X is like*

$$\mathcal{B} = \{ \{x \in X : p_{i_1}(x) < \varepsilon, \dots, p_{i_n}(x) < \varepsilon\} : i_1, \dots, i_n \in I, n \in \mathbb{N}, \varepsilon > 0 \}$$

Conversely, the topology of any locally convex t.v.s. can be generated by a family of seminorms.

Proof. Each element in the subbasis of \mathcal{T}_P is like $\{x \in X : p_i(x) < \varepsilon\} = \varepsilon \mathring{U}_{p_i}$, which is clearly absorbing and absolutely convex. Therefore, every element in \mathcal{B} is convex. (X, \mathcal{T}_P) is a locally convex t.v.s..

Conversely, if (X, \mathcal{T}) is a locally convex t.v.s., the basis of neighbourhoods of the origin in X consists of absorbing and absolutely convex sets, which can generate a family of seminorms by above proposition. Therefore, these seminorms can generate a locally convex topology \mathcal{T}_P . In fact, $\mathcal{T}_P = \mathcal{T}$, since $\mathring{U}_{p_A} \subset A \subset U_{p_A}$. \square

Remark. By this theorem, in a vector space X , the seminorms on X can coincide with a locally convex topology making X be a t.v.s..

There some extra properties for the seminorms on a vector space.

Proposition 5. *Let X be a vector space and p be a seminorm on X . Then,*

$$1) \forall r > 0, r \mathring{U}_p = \{x \in X : p(x) < r\} = \mathring{U}_{\frac{1}{r}p}.$$

$$2) \forall x \in X, x + \mathring{U}_p = \{y \in X : p(y - x) < 1\}.$$

3) if q is a seminorm on X , $p \leq q \Leftrightarrow \mathring{U}_q \subset \mathring{U}_p$.

4) if $\{s_i\}_{i=1}^n$ are seminorms on X , then $s(x) = \max_{i=1, \dots, n} s_i(x)$ is also a seminorm and $\mathring{U}_s = \bigcap_{i=1}^n \mathring{U}_{s_i}$

Theorem 6. Let $\mathcal{P} = \{p_i\}_{i \in I}$ and $\mathcal{Q} = \{q_j\}_{j \in J}$ be two families of seminorms on a vector space X inducing \mathcal{T}_P and \mathcal{T}_Q , then

$$\mathcal{T}_Q \subset \mathcal{T}_P \Leftrightarrow \forall q \in \mathcal{Q}, \exists \{i_k\}_{k=1}^n \subset I, \exists C > 0, \text{ s.t. } Cq(x) \leq \max_{k=1, \dots, n} p_{i_k}(x)$$

Proof. This right side of above statement is equivalent to that

$$\forall q \in \mathcal{Q}, \exists \{i_k\}_{k=1}^n \subset I, \exists C > 0, \text{ s.t. } C \bigcap_{k=1}^n \mathring{U}_{p_{i_k}} \subset \mathring{U}_q$$

So, it is clearly equivalent to $\mathcal{T}_Q \subset \mathcal{T}_P$. □

Remark. Because of this, we have the definition of equivalent norms. In fact, two norms p and q are said to be equivalent if and only if there exists $C_1, C_2 > 0$, s.t. $C_1 p(x) \leq q(x) \leq C_2 p(x)$, for all $x \in X$. This definition means that two equivalent norms can generate one same topology.

Corollary 7. The family $\mathcal{P} = \{p_i\}_{i \in I}$ of seminorms and $\mathcal{Q} = \{ \max_{i \in B} p_i : \emptyset \neq B \subset I \text{ with } B \text{ finite} \}$ can generate one same topology on X

2.2.3 Separability and Metrizable

For a t.v.s., (T_1) is equivalent to Hausdorff and (T_1) is associated with the ability of a topology separating points. Therefore, wheather a t.v.s. is a Hausdorff space or not is completely determined by wheather the topology of it can separate points or not. Then for a locally convex t.v.s., whose topology is induced by a family of seminorms, this separability is related to these seminorms.

Definition 1. A family of seminorms $\mathcal{P} = \{p_i\}_{i \in I}$ on a vector space X is said to be separating, if

$$\forall x \in X \setminus \{0\}, \exists i \in I \text{ s.t. } p_i(x) \neq 0$$

Remark. In fact, above condition is equivalent to

$$\text{if } p_i(x) = 0, \forall i \in I \Rightarrow x = 0$$

Now, we can give the condition that makes a locally convex t.v.s. be Hausdorff.

Theorem 2. *A locally convex t.v.s. X is Hausdorff if and only if its topology can be induced by a separating family of seminorms $\mathcal{P} = \{ p_i \}_{i \in I}$.*

Proof. If $\mathcal{P} = \{ p_i \}_{i \in I}$ is separating, the fact that \mathcal{T}_P is Hausdorff can be obtained easily by the definition.

Conversely, if X is Hausdorff, for $x \neq 0$, we can find a $U \in \mathcal{F}(0)$, s.t. U can separate x and 0. But since X is locally convex, U can be chosen as \mathring{U}_p for a seminorm p . Thus, for this p , $p(x) \neq 0$. \square

For the metrizable of a locally convex t.v.s., the consequence is also easier than general case.

Theorem 3. *A locally convex t.v.s. X is metrizable if and only if its topology is determined by a countable separating family of seminorms.*

Proof. It can be directly obtained by the Nagata–Smirnov’s Metrization Theorem. Also, there is a more explicit proof. If the topology of X is generated by a countable separating family of seminorms $\mathcal{P} = \{ p_n \}_{n=1}^\infty$, we can define the metric d on X by

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

Conversely, if (X, d) is the metric space, the subbasis of the topology generated by this metric is like $U_n = \{x: d(x, 0) < 1/n\}$. And these U_n can provide a countable separating family of seminorms. In fact, if $\mathcal{Q} = \{ q_i \}_{i \in I}$ generates the topology of X , for each U_n , there are $q_1, \dots, q_k \in \mathcal{Q}$ and $\varepsilon_1, \dots, \varepsilon_k > 0$, s.t. $\bigcap_{i=1}^k \{x: q_i(x) < \varepsilon_i\} \subset U_n$. Then let $p_n = \sum_{i=1}^k \varepsilon_i^{-1} q_i$. It can check that the family $\{p_n\}_{n=1}^\infty$ generate the coincided topology on X . \square

2.2.4 Continuous Linear Maps on LCTVS

To give the special property of continuous linear maps on LCTVS, we need firstly refine the family of seminorms.

Definition 1. A family $\mathcal{Q} = \{ q_j \}_{j \in J}$ of seminorms on a vector space X is said to be directed if

$$\forall j_1, \dots, j_n \in \mathcal{Q}, \exists j \in J \text{ \& } C > 0, \text{ s.t. } Cq_j(x) \geq \max_{k=1, \dots, n} q_{j_k}(x), \forall x \in X$$

Remark. By **Proposition 5** in **1.2.2**, this definition is equivalent to that

$$\forall \mathring{U}_{q_{j_1}}, \dots, \mathring{U}_{q_{j_n}}, \exists \mathring{U}_{q_j} \text{ s.t. } \mathring{U}_{q_j} \subset \bigcap_{k=1}^n \mathring{U}_{q_{j_k}}$$

And thus the basis of this directed family of seminorms should be like

$$\mathcal{B}_d = \{ r\mathring{U}_q: q \in \mathcal{Q}, r > 0 \}$$

By this special topology, we can find the condition making linear functional continuous.

Proposition 2. *Let \mathcal{T} be a locally convex topology on a vector space X generated by a directed family \mathcal{Q} of seminorms on X . Then*

$$L: X \longrightarrow \mathbb{K}$$

is a \mathcal{T} -continuous linear functional if and only if $\exists q \in \mathcal{Q}$ s.t. L is q -continuous, i.e. $|L(x)| \leq Cq(x)$ for some $C > 0$.

Proof. In fact, this property of continuous linear functional is because the element in a directed locally convex topology is like $r\mathring{U}_q$. In fact, we just need to check the origin point.

If L is continuous, there exists a $r\mathring{U}_q$ s.t. $r\mathring{U}_q \subset L^{-1}(B_1(0))$, where $B_1(0)$ is the unit ball centered at 0. This is equivalent to the q -continuity of L .

Conversely, it is clearly by the fact $\mathcal{T}_q \subset \mathcal{T}$. \square

We can easily see that the topology of a locally convex t.v.s. can be always induced by a directed family of seminorms by the **Corollary 7** in **1.2.2**. Thus, we have the corollary.

Corollary 3. *(X, \mathcal{T}) is a locally convex t.v.s. and \mathcal{T} is generated by the family $\mathcal{P} = \{p_i\}_{i \in I}$. Then $L: X \longrightarrow \mathbb{K}$ is a continuous linear functional if and only if*

$$\exists i_1, \dots, i_n \in I, \exists C > 0 \text{ s.t. } |L(x)| \leq C \max_{k=1, \dots, n} p_{i_k}(x), \forall x \in X$$

And this corollary can be easily extended to linear maps. And the proof is similar as above statement. In fact, we just need to replace $B_1(0)$ by \mathring{U}_q .

Theorem 4. *Let X and Y be two locally convex t.v.s.'s generated by \mathcal{P} and \mathcal{Q} . Then linear map $f: X \rightarrow Y$ is continuous if and only if*

$$\forall q \in \mathcal{Q}, \exists p_1, \dots, p_n \in \mathcal{P}, \exists C > 0 \text{ s.t. } q(f(x)) \leq C \max_{i=1, \dots, n} p_i(x)$$

2.3 The Hahn-Banach Theorem

2.3.1 Two Forms of Hahn-Banach Theorem

Theorem 1 (Geometric form). *Let X be a t.v.s. over \mathbb{K} , N be a linear subspace of X and Ω be a convex open subset of X with $N \cap \Omega = \emptyset$. Then there is a closed hyperplane H of X s.t.*

$$N \subset H \quad \& \quad H \cap \Omega = \emptyset$$

Proof. Assume that $\Omega \neq \emptyset$.

Let $\mathcal{F} = \{ \text{all linear subspace } L \text{ of } X \text{ s.t. } N \subset L \text{ and } L \cap \Omega = \emptyset \}$.

Since $N \in \mathcal{F}$, $\mathcal{F} \neq \emptyset$. And \mathcal{F} can be ordered by " \subset ". Clearly, for every chain $\mathcal{C} = \{C_i\}_{i \in I} \in \mathcal{F}$, it has a maximal element $C = \bigcup_{i \in I} C_i \in \mathcal{F}$. Then by Zorn's Lemma, there exists a maximal $H \in \mathcal{F}$. And this H is closed because the maximality of H .

In fact, we just need to check $\dim X/H = 1$. In the case $\mathbb{K} = \mathbb{R}$, if $\dim X/H \geq 2$, then we can find one dimensional subspace L in X s.t. $L \oplus H$ satisfies above condition. Then it is a contradiction to the maximality of H . And for $\mathbb{K} = \mathbb{C}$, using above process to get a real hyperplane H_0 and then $H = H_0 \cap iH_0$ is the hyperplane we need.

The main difficulty of this proof is to construct L . In fact, we just need to find a line \tilde{L} in X/H s.t. $\tilde{L} \cap A = \emptyset$, where $A = \bigcup_{\lambda > 0} \lambda \pi(\Omega)$ is a cone. This \tilde{L} can easily be found because we can find $\bar{x}, -\bar{x} \in X/H \setminus A$ with $\bar{x} \neq \bar{0}$ by using the fact that $\dim X/H \geq 2$. \square

Remark. In this theorem, if N is an affine linear subspace, then the H can be chosen as a affine hyperplane satisfying above condition. It is clear by using translation.

Theorem 2 (Analytic form). *Let p be a seminorm on a vector space X over \mathbb{K} , M is a linear space of X and f is a linear functional on M s.t. $|f(x)| \leq p(x)$, $\forall x \in M$. Then there is a linear function \tilde{f} on X s.t. $\tilde{f}(x) = f(x)$, $\forall x \in M$ and $|\tilde{f}(x)| \leq p(x)$, $\forall x \in X$.*

Proof. Let $N = \{x \in M : f(x) = 1\}$ (affine hyperplane) and $\Omega = \{x \in X : p(x) < 1\}$ (open convex set). Then since $N \cap \Omega = \emptyset$, we can find a closed affine hyperplane H of X , s.t. $N \subset H$ and $H \cap \Omega = \emptyset$. Fixed $x_0 \in N \subset H$, $H - x_0$ is a hyperplane, thus it can generate a functional \tilde{f} on X . Set $\tilde{f}(x_0) = 1$, then define

$$\begin{aligned} \tilde{f} : \quad X = (H - x_0) \oplus \mathbb{K}x_0 &\longrightarrow \mathbb{K} \\ (h - x_0) + \lambda x_0 &\longmapsto \lambda \end{aligned}$$

Since $M = (N - x_0) \oplus \mathbb{K}x_0$ and $f(x_0) = 1$ and $H \cap \Omega = \emptyset$, we can easily check that $\tilde{f}(x) = f(x)$, $\forall x \in M$ and $|\tilde{f}(x)| \leq p(x)$, $\forall x \in X$. \square

2.3.2 Applications

By the geometric form of the Hahn-Banach Theorem, it says about hyperplanes can separate some non-intersecting subsets. And moreover, by the analytic form of the Hahn-Banach Theorem, since hyperplanes are linked with functionals, functionals can also separate some non-intersecting subsets in a t.v.s. First, we can define this association more explicit.

Definition 1. Let X be a t.v.s. over \mathbb{R} and H is a closed affine hyperplane. For $A, B \subset X$ and $A \cap B = \emptyset$, we say A and B separated by H , if

$$\exists a \in \mathbb{R}, \text{ s.t. } H = f^{-1}(\{a\}) \text{ for some } f: X \rightarrow \mathbb{R} \text{ \& } f(A) \geq a, f(B) \leq a.$$

Theorem 2. Let X be a t.v.s. over \mathbb{R} and A, B be two disjoint convex non-empty subsets of X , then we have:

- 1) if A is open, then there exists a closed affine hyperplane H of X separating A and B .
- 2) if A, B are open, H can strictly separate A and B .
- 3) if A is a cone and B is open, then H can be chosen as a hyperplane.

Proof. Let $U = A - B$. Clearly, U is open and convex and $N = \{0\} \cap H = \emptyset$. Therefore, there exists a continuous functional f on X , s.t. $f(U) > 0$ i.e. $f(x) > f(y), \forall x \in A \text{ \& } y \in B$. Since $B \neq \emptyset$, let $a = \inf_{x \in A} f(x) > -\infty$. And thus 1) and 2) can be obtained.

For 3), if A is a cone, then $tf(x) = f(tx) \geq a, \forall t > 0 \Rightarrow f(x) \geq \frac{a}{t}$. Thus, $f(A) \geq 0$. H can be chosen as a hyperplane. \square

For a locally convex t.v.s X , each point in X has some convex neighbourhoods. Therefore, we have following corollaries.

Corollary 3. Let X be a locally convex t.v.s over \mathbb{R} .

- 1) If A and B are two disjoint closed subsets and B is compact, then A and B are strictly separated.
- 2) If A is a closed convex subset of X and $x \notin A$, then x is strictly separated from A .
- 3) If A is a subset of X , then $\overline{\text{span}\{A\}}$ is the intersection of all closed hyperplanes containing A .

And these consequences can be extended to \mathbb{C} .

Theorem 4. Let X be a locally convex t.v.s over \mathbb{C} and A, B be two disjoint closed convex subsets of X . If B is compact, then there is a continuous linear functional $f: X \rightarrow \mathbb{C}$, and $\alpha \in \mathbb{R}$, and $\varepsilon > 0$ s.t.

$$\text{Re } f(x) \leq \alpha < \alpha + \varepsilon \leq \text{Re } f(y), \forall x \in A, \forall y \in B$$

2.4 Banach Spaces

The Banach space is a very special locally convex topological space, whose topology is generated by only one seminorm. To make this space be a Hausdorff space, this seminorm is actually a norm. And more, we need it become complete. So, a Banach space is a complete normed space. Because it is definitely a locally convex Hausdorff t.v.s., all results mentioned above can be applied on it. And we can have more interesting results of the Banach space because of its simple structure.

2.4.1 Elementary Properties

Definition 1. A normed space is a vector space X with a compatible norm $\|\cdot\|$, which makes it be a locally convex Hausdorff topological vector space. A Banach space is a complete normed space.

Remark. By this definition, we know that any result mentioned above can be also true for the Banach space $(X, \|\cdot\|)$ by replacing the family of seminorms by the $\|\cdot\|$. Moreover, two equivalent norms on X provide same topology.

The properties of finite dimensional Banach spaces is in the subsection 1.1.5. And for the quotient space of a Banach space, we can get a more explicit expression of the induced quotient norm.

Theorem 2. Let $(X, \|\cdot\|)$ be a Banach space and M be a closed linear subspace of X and $\pi: X \rightarrow X/M$, then the induced quotient norm on X/M is defined as $\|x + M\| = \inf \{\|x + y\| : y \in M\}$. Thus, using the consequences of the quotient t.v.s., we have following results.

- 1) π is continuous and $\|\pi(x)\| \leq \|x\|$.
- 2) X/M is a Banach space.
- 3) $W \subset X/M$ is open if and only if $\pi^{-1}(W)$ is open.
- 4) π is open.
- 5) if N is a finite dimensional subspace of X , then $M + N$ is closed.

Proof. The element in the subbasis of the topology generated by $\|\cdot\|$ at 0 is like $U_\varepsilon = \{x \in X : \|x\| < \varepsilon\}$ for a fixed $\varepsilon > 0$. Therefore, the element V_ϵ in the subbasis of the induced topology on X/M at 0 satisfies $\pi^{-1}(V_\epsilon) = \{x : \|x\| < \epsilon\}$ for some $\epsilon > 0$, i.e. $V_\epsilon = \{x + M : \|x + y_0\| < \epsilon, \exists y_0 \in M\}$. Clearly, V_ϵ is an absorbing absolutely convex set. Then the Minkowski functional p_ϵ of V_ϵ is like

$$\begin{aligned} p_\epsilon(x + M) &= \inf\{\lambda > 0 : x + M \in \lambda V_\epsilon\} \\ &= \inf\{\lambda > 0 : \|x + y_0\| < \lambda\epsilon, \exists y_0 \in M\} \\ &= \frac{1}{\epsilon} \inf\{\|x + y\| : y \in M\} \end{aligned}$$

Clearly, for any $\epsilon > 0$, p_ϵ is a norm and it is equivalent to the norm $\|x + M\|$. Thus the quotient topology is definitely generated by $\|x + M\|$. Then, 1) and 3) clearly hold by the definition. 2) is true since H is a Banach space. 4) is true for any general t.v.s.. 5) holds since $N + M$ is a finite dimensional subspace of X/M . \square

Similarly, we can define the norm on the product space of some Banach spaces.

Definition 3. Let $\{X_i\}_{i=1}^p$ is a family of Banach spaces with norm $\|\cdot\|$. Then,

1) if $1 \leq p < \infty$,

$$\oplus_{i=1}^p X_i = \left\{ x \in \prod_i^p X_i : \|x\| = \left[\sum_{i=1}^p \|x_i\|^p \right]^{1/p} \right\}$$

2) if $p = \infty$,

$$\oplus_{i=1}^\infty X_i = \left\{ x \in \prod_i^p X_i : \|x\| = \sup_i \{x(i)\} < \infty \right\}$$

2.4.2 Linear Transformations and Linear Functionals

By results in **Theorem 4** in the subsection **1.2.4** the linear transformation $T: X \rightarrow Y$ between two Banach spaces is continuous if and only if $\exists C > 0$, s.t. $\|Tx\| \leq \|x\|$, $\forall x \in X$. Because of this property, we can define the norm of continuous linear transformation.

Definition 1. Let $T: X \rightarrow Y$ between two Banach spaces be a linear transformation. Then the norm of T is defined as

$$\|T\| = \sup_{x \in X, x \neq 0} \frac{\|Tx\|}{\|x\|}$$

Remark. By this definition, we know T is continuous if and only if $\|T\| < \infty$, and we call T is bounded. Also, the formula of the norm can be

$$\|T\| = \sup_{\|x\|=1} \frac{\|Tx\|}{\|x\|}$$

Let $\mathcal{B}(X, Y) = \{\text{all bounded linear transformations between } X \text{ and } Y\}$. $\mathcal{B}(X, Y)$ is a Banach space with this norm. Then in the space $\mathcal{B}(X, Y)$, there are two topologies. The first one is generated by a family of seminorms $\{p_x\}_{x \in X}$, where $p_x(T) = \|Tx\|$. And the second one is generated by the norm $\|\cdot\|$. In fact, the convergence of a sequense in $\mathcal{B}(X, Y)$ with respect to the first one is about pointwise convergence, and with respect to the second one is about uniform convergence. Clearly, the second one is stronger. But in some case the first one can give some information of the first one.

Theorem 2 (The Principle of Uniform Boundedness). *Let X be a Banach space and Y be a normed space, and $\{T_i\}_{i \in I}$ be a subset of $\mathcal{B}(X, Y)$. If for any $x \in X$ $\{\|T_i x\|\}_{i \in I}$ is bounded, then $\{\|T_i\|\}_{i \in I}$ is bounded.*

Proof. Firstly, there is a open ball $B(x_0, \varepsilon)$ such that for any $x \in B(x_0, \varepsilon)$, $\|T_i\|x \leq K$ for some constant K . Assume it is not true. We can find a family of open balls $B(x_n, \varepsilon_n)$ s.t. $B(x_n, \varepsilon_n) \subset B(x_{n-1}, \varepsilon_{n-1})$, $\varepsilon_n < 1/n$, and a sequence $i_n \subset I$, satisfying $\|T_{i_n} x\| > n$. And since each $\overline{B(x_n, \varepsilon_n)}$ is compact, there is a $z \in \bigcap_{n=1}^{\infty} \overline{B(x_n, \varepsilon_n)}$. This contradicts to the boundedness of $\{T_i z\}_{i \in I}$.

Next, with this $B(x_0, \varepsilon)$, if $y \in X$ and $y \neq 0$, then $z = \frac{\varepsilon}{\|y\|}y + x_0 \in B(x_0, \varepsilon)$. Therefore,

$$\frac{\varepsilon}{\|y\|} \|T_i y\| - \|T_i x_0\| \leq \left\| \frac{\varepsilon}{\|y\|} \|T_i y\| + \|T_i x_0\| \right\| = \|T_i z\| \leq K$$

Then since $K' = \sup_{i \in I} \|T_i x_0\| < \infty$,

$$\|T_i y\| \leq \frac{K + K'}{\varepsilon} \|y\| < \infty$$

Thus $\|T_i\|_{i \in I}$ is bounded. □

Remark. Let $X^* = \mathcal{B}(X, \mathbb{K})$ be the continuous linear functional space. Above theorem can easily apply to it.

Corollary 3. *For a Banach space X and $S \subset X^*$ a subset, S is normed bounded if and only if for every $x \in X$, $\sup\{|f(x)| : f \in S\} < \infty$.*

Another important property of the linear transformation between two Banach spaces is that any continuous linear surjective transformation is open.

Theorem 4 (Open Mapping Theorem). *Let X and Y be two Banach spaces and $T: X \rightarrow Y$ be a continuous linear surjection. Then T is open.*

Proof. Let X_ε and Y_ε be two open balls in X and Y with centering at 0 and radius ε .

1) For any $\varepsilon > 0$, there is a $\delta > 0$, s.t. $Y_\delta \subset \overline{T(X_\varepsilon)}$.

By $X = \bigcup_{n=1}^{\infty} nX_\varepsilon$, $Y = \bigcup_{n=1}^{\infty} T(nX_\varepsilon)$. Then since Y is the second Baire category, there exists n_0 and a ball $B_r(z) \subset Y$ s.t. $B_r(z) \subset \overline{T(n_0 X_\varepsilon)}$ i.e. $B_\delta(y_0) \subset \overline{T(X_\varepsilon)}$ with $\delta = r/n$ and $y_0 = z/n$. Then since for any $y \in Y_\delta$, $y = (y + y_0) - y_0$

$$Y_\delta \subset \{y_1 - y_2 : y_1, y_2 \in B_\delta(y_0)\} \subset \overline{T(\{x_1 - x_2 : x_1, x_2 \in X_\varepsilon\})} \subset \overline{T(X_{2\varepsilon})}$$

2) For any $\varepsilon_0 > 0$, there is a $\delta_0 > 0$, s.t. $Y_{\delta_0} \subset T(X_{2\varepsilon_0})$.

Choose a positive sequence $\{\varepsilon_n\}$ with $\sum_{n=1}^{\infty} \varepsilon_n < \varepsilon_0$, then there exists a positive sequence δ_n with $\delta_n \rightarrow 0$ s.t. $Y_{\delta_n} \subset \overline{T(X_{\varepsilon_n})}$

For any $y \in Y_{\delta_0}$, there is a $x_0 \in X_{\varepsilon_0}$, s.t. $\|y - Tx_0\| < \delta_1$. since $y - Tx_0 \in Y_{\delta_0}$, there is a $x_1 \in X_{\varepsilon_1}$ s.t. $\|y - Tx_0 - Tx_1\| < \delta_2$. Thus by induction, we can find a sequence $\{x_n\}$ s.t. $x_n \in X_{\varepsilon_n}$ and

$$\left\| y - T \left(\sum_{k=0}^n x_k \right) \right\| < \delta_{n+1}$$

And since $\sum_{n=1}^{\infty} \varepsilon_n < \varepsilon_0$, the sequence $\sum_{k=0}^n x_k$ absolutely converges to x in X . And

$$\|x\| \leq \sum_{n=1}^{\infty} \|x_n\| \leq \sum_{n=1}^{\infty} \varepsilon_n < 2\varepsilon_0$$

By the continuity of T and $\delta_n \rightarrow 0$, $y = Tx$. Therefore, $Y_{\delta_0} \subset T(X_{2\varepsilon_0})$.

3) For any open set $G \subset X$ and $y_0 = Tx_0$ with $x_0 \in G$, there is a open ball Y_δ s.t. $y_0 + Y_\delta \subset T(G)$.

Since G is open, there is a open ball X_ε with $x_0 + X_\varepsilon \subset G$. Then there is a $\delta > 0$ s.t. $Y_\delta \subset T(X_\varepsilon)$. Thus

$$y_0 + Y_\delta \subset Tx_0 + T(X_\varepsilon) = T(x_0 + X_\varepsilon) \subset T(G) \quad \square$$

Then there is a directed corollary from the Open Mapping Theorem.

Corollary 5 (The Inverse Mapping Theorem). *Let X and Y be two Banach spaces and $T: X \rightarrow Y$ be a continuous linear isomorphism, then T^{-1} is also bounded.*

Corollary 6 (The Closed Graph Theorem). *Let X and Y be two Banach spaces and $T: X \rightarrow Y$ be a linear transformation. If the graph of T , $G = \{x \oplus Tx: x \in X\}$ is closed, then T is continuous.*

Proof. Since X and Y are two Banach spaces and G is closed, G is also a Banach space. Let $P_1: G \rightarrow X$ be $P_1(x \oplus Tx) = x$ and $P_2: G \rightarrow Y$ be $P_2(x \oplus Tx) = Tx$. Then P_1 is bounded and bijective. Therefore, by the Inverse Mapping Theorem, $P_1^{-1}: X \rightarrow G$ is continuous. Since P_2 is also continuous, $A = P_2 \circ P_1^{-1}$ is continuous. \square

The next important theorem is Hahn-Banach Theorem. We have known the Hahn-Banach Theorem holds in any general t.v.s.. But for the Banach space, it can induce some interesting corollaries of the bounded linear functional space.

Corollary 7. *Let X be a Banach space. Then,*

- 1) if $\{x_i\}_{i=1}^d$ is a linearly independent in X , and $\{\alpha_i\}_{i=1}^d$ are arbitrary scalars, then there is a $f \in X^*$ s.t. $f(x_i) = \alpha_i$ for $i = 1, \dots, d$.
- 2) if Y is a linear subspace of X and $x_0 \in X$ with $\inf\{\|x_0 - y\| : y \in Y\} = d > 0$, then there is a $f \in X^*$ s.t. $f(x_0) = 1$ and $f(y) = 0 \forall y \in Y$ and $\|f\| = d^{-1}$.
- 3) if $x \in X$, then

$$\|x\| = \sup_{\|f\| \neq 0} \frac{|f(x)|}{\|f\|} = \sup\{|f(x)| : f \in X^*, \|f\| \leq 1\}$$

- 4) if $x \neq y$ in X , then there is a $f \in X^*$ s.t. $f(x) \neq f(y)$.
- 5) if Y is a linear subspace of X and Y is not dense, then there is $f \in X^*$ with $f \neq 0$ s.t. $f(y) = 0 \forall y \in Y$.
- 6) if Y is a linear subspace of X , then

$$\overline{Y} = \bigcap \{\ker f : f \in X^*, Y \subset \ker f\}$$

Proof. For 1), let $Y = \text{span}\{x_i\}_{i=1}^d$ and g be a linear functional on Y with $g(x_i) = \alpha_i$, then g can extend to f .

For 2), let $Y_1 = \text{span}\{Y, \{x_0\}\}$ and g be a linear functional on Y_1 with $g(y + \lambda x_0) = \lambda$.

Then since $\|y + \lambda x_0\| = |\lambda| \|\frac{1}{\lambda}y + x_0\| \geq |\lambda| d$, $\|g\| \leq d^{-1}$. Let $\{y_n\} \in Y$ s.t. $\|x_0 - y_n\| \rightarrow d$, then $1 = g(x_0 - y_n) \leq \|g\| \|x_0 - y_n\| \rightarrow \|g\| d$.

Thus $\|g\| = d^{-1}$. g can extend to f on X .

For 3), clearly, $\|x\| \geq \sup\{|f(x)| : f \in X^*, \|f\| \leq 1\}$. But by above proof, there is a $f(x) = \|x\|$ with $\|f\| = 1$.

4) and 5) can be easily obtained by the separability of continuous linear functionals and the fact that X is locally convex.

6) is a special case of **Corollary 3** in the subsection **1.3.2**. \square

For linear functionals on a general vector space, there is a interesting result.

Proposition 8. Let $f, \{f_k\}_{k=1}^n$ be linear functionals on a vector space X . If $\bigcap_{k=1}^n \ker f_k \subset \ker f$, then there are scalars $\alpha_1, \dots, \alpha_n$ such that $f = \sum_{k=1}^n \alpha_k f_k$.

Proof. Assume for $1 \leq k \leq n$, $\bigcap_{j \neq k} \ker f_j \neq \bigcap_{j=1}^n \ker f_j$.

Therefore, there is a $y_k \in \bigcap_{j \neq k} \ker f_j$ s.t. $y_k \notin \bigcap_{j=1}^n \ker f_j$, i.e. $f_j(y_k) = 0$ for $j \neq k$ and $f_k(y_k) \neq 0$.

Thus we can find $x_{k=1}^n$, s.t. $f_k(x_k) = 1$ and $f_j(x_k) = 0$ for $j \neq k$.

Let $\alpha_k = f(x_k)$. For $x \in X$, let $y = x - \sum_{k=1}^n \alpha_k f_k(x) x_k$.

Then $f_j(y) = f_j(x) - \sum_{k=1}^n \alpha_k f_k(x) f_j(x_k) = 0$. Thus $f(y) = 0$,

i.e. $f(x) = \sum_{k=1}^n \alpha_k f_k(x)$. \square

2.4.3 Weak Topologies

In a Banach space X , X^* denoted the set of all bounded linear functionals is also a Banach space with the induced norm. Then we can also find all bounded linear functionals on the space $(X^*, \|\cdot\|)$, denoted by X^{**} . By the Hahn-Banach Theorem, we can easily know the following proposition.

Proposition 1. *Let X be a Banach space. Then X can be isometrically embeded in X^{**} .*

Proof. Define the map ϕ , where $\hat{x}(f) = f(x)$,

$$\begin{aligned} \phi: X &\longrightarrow X^{**} \\ x &\longmapsto \hat{x} = \phi(x) \end{aligned}$$

Then by the corollary of Hahn-Banach Theorem, $\|\hat{x}\| = \|x\|$. \square

Remark. In a special case, if $X^{**} = \hat{X}$, where $\hat{X} = \pi(X)$, then X is called a reflexive space.

For a Banach space X and X^* , we can define different topologies.

Definition 2. Let X be a Banach space.

- 1) The weak topology on X , denoted by wk , is generated by the family of seminorms $\{p_f: f \in X^*\}$, where $p_f(x) = |f(x)|$.
- 2) The weak* topology on X^* , denoted by wk^* , is generated by the family of seminorms $\{p_x: x \in X\}$, where $p_x(f) = |f(x)|$.

Remark. Since the $|f(x)| \leq \|f\| \|x\|$, we can easily know wk on X and wk^* on X^* are weaker than the norm topology respectively.

- 1) The subbasis of the weak topology at x_0 is like

$$U_\varepsilon(x_0) = \{x \in X: |f(x - x_0)| < \varepsilon\}$$

Therefore, a net $\{x_\alpha\}$ in X converges weakly to x_0 if and only if $f(x_\alpha) \rightarrow f(x_0)$ for all $f \in X^*$.

- 2) The subbasis of the weak* topology at f_0 is like

$$V_\varepsilon(f_0) = \{f \in X^*: |(f - f_0)(x)| < \varepsilon\}$$

Therefore, a net $\{f_\alpha\}$ in X^* converges weakly to f_0 if and only if $f_\alpha(x) \rightarrow f_0(x)$ for all $x \in X$.

There are some easy properties between these two topologies and the respective norm topology.

Proposition 3. *Let X be a Banach space. Then we have*

$$1) (X, wk)^* = X^*$$

$$2) (X^*, wk^*)^* = X$$

$$3) \text{ if } A \subset X \text{ is convex, then } \overline{A} = \overline{A}^{wk}$$

Proof. 1) can be obtained easily by the fact that wk is weaker than the norm topology and by the definition of continuous functionals.

For 2), because of the **Corollary 3** in the subsection **1.2.4** and the **Proposition 5** in the subsection **1.2.2** and the **Proposition 8**.

$$\begin{aligned} F \in (X^*, wk^*)^* &\Rightarrow |F(f)| \leq \sum_{i=1}^n |\hat{x}_k| \text{ for some } \{\hat{x}_k\}_{k=1}^n \subset \hat{X} \\ &\Rightarrow \bigcap_{k=1}^n \ker \hat{x}_k \subset \ker F \Rightarrow F = \sum_{k=1}^n \alpha_k \hat{x}_k \in \hat{X} \cong X \end{aligned}$$

And the converse is because wk^* is weaker than norm topology.

For 3), because wk is weaker than the norm topology, $\overline{A} \subset \overline{A}^{wk}$. By the separability of continuous functionals in **Theorem 4** in the **1.3.2**, we can find a $f \in X^*$ s.t. f separates \overline{A} and any $x \in X \setminus \overline{A}$ in using a positive real number α . But $\overline{A} \subset B = \{y \in X : \operatorname{Re} f(y) \leq \alpha\}$ and B is clearly wk -closed. Therefore, $\overline{A}^{wk} \subset B$. $x \notin \overline{A}$ implies $x \notin B$, thus $x \notin \overline{A}^{wk}$, i.e. $\overline{A}^{wk} \subset \overline{A}$. \square

Also, by the Hahn-Banach Theorem and the **Theorem 4** in the **1.2.4**, the bounded linear map can be characterized as the following statement.

Proposition 4. *Let $T: X \rightarrow Y$ be a linear map between two Banach spaces. Then T is bounded if and only if $T: (X, wk) \rightarrow (Y, wk)$ is continuous.*

By the **Theorem 3** in the **1.2.3**, we can easily see how to make X^* be wk^* -metrizable.

Theorem 5. *If X is a Banach space, then X^* is wk^* -metrizable if and only if X is separable.*

Proof. It is because that wk^* is generated by $\{p_x : x \in X\}$. \square

Let Y be a closed subspace a Banach space X . We can define the orthogonal complement of Y , $Y^\perp = \{f \in X^* : Y \subset \ker f\}$ (similar definition in X^*). Then we have the following theorem.

Theorem 6. *If Y is a closed subspace a Banach space X and $\pi: X \rightarrow X/Y$ is the quotient map, then*

1) *the following map is an isometric isomorphism, i.e. $X^*/Y^\perp \cong Y^*$.*

$$\begin{aligned} \rho: \quad X^*/Y^\perp &\longrightarrow Y^* \\ f + Y^\perp &\longmapsto f|_Y \end{aligned}$$

2) the the following map is an isometric isomorphism, i.e. $(X/Y)^* \cong Y^\perp$.

$$\begin{aligned} \kappa: (X/Y)^* &\longrightarrow Y^\perp \\ f &\longmapsto f \circ \pi \end{aligned}$$

Proof. For 1), clearly, ρ is linear and injective. Since for any $f \in X^*$ and $g \in Y^\perp$,

$$\|f|_Y\| = \|(f + g)|_Y\| \leq \|(f + g)\|$$

we have $\|f|_Y\| \leq \|f + Y^\perp\|$. By the Hahn-Banach Theorem, for any $\phi \in Y^*$, there is a $f \in X^*$ with $f|_Y = \phi$ and $\|\phi\| = \|f\| \geq \|f + Y^\perp\|$.

For 2), clearly $\|\kappa(f)\| \leq \|f\|$. We can find a sequence $x_n + y_n$ with $x_n \in X$ and $y_n \in Y$ and $\|x_n + y_n\| < 1$ s.t.

$$\|\kappa(f)\| \geq |\kappa(f)(x_n + y_n)| = f(x_n + Y) \rightarrow \|f\|$$

thus κ is an isometry. And by the universal property of quotient map, in the **Proposition 1** in the subsection **1.1.3**, we can prove that κ is surjective.

$$\begin{array}{ccc} X & \xrightarrow{f} & \mathbb{K} \\ \downarrow \pi & \nearrow \tilde{f} & \uparrow f' \\ X/\ker f & \xrightarrow{i} & X/Y \end{array}$$

□

We have an important theorem to describe the wk^* -compactness in dual space.

Theorem 7 (Alaoglu's Theorem). *Let X be a normed space, then the unit ball in X^* is wk^* -compact.*

Proof. Let $D_x = \{\alpha \in \mathbb{K} : |\alpha| \leq 1\}$ for any $x \in X$, then put $D = \prod_{x \in X} D_x$. By Tychonoff's Theorem, D is compact. Let the unit ball in X^* be denoted by B and τ be defined as

$$\begin{aligned} \tau: B &\longrightarrow D \\ f &\longmapsto \tau(f) \end{aligned}$$

where $\tau(f)(x) = f(x)$. Then we can prove that B is homeomorphic to $\tau(B)$ with respect to the induced topology of D and $\tau(B)$ is closed in D . Thus B is compact.

1) Injection: $\tau(f) = \tau(g) \Rightarrow f(x) = g(x) \forall x \in B \Rightarrow f = g$.

2) Continuity: $f_i \rightarrow f \text{ } wk^* \Rightarrow \tau(f_i)(x) \rightarrow \tau(f)(x) \forall x \in B \Rightarrow \tau(f_i) \rightarrow \tau(f)$.

- 3) Closedness: $\tau(f_i) \rightarrow f \in D \Rightarrow f_i(x) \rightarrow f(x) \forall x \in B$, then we extend f by defining $\tilde{f}(x) = \alpha^{-1}f(\alpha x)$ for any $x \in X$ and $\alpha > 0$ with $\|\alpha x\| \leq 1$. It is well-defined by the linearity. Then $\tilde{f} \in B$. $\tau(B)$ is closed, and $\tau(B)$ is complete.
- 4) Homeomorphism: $\tau: B \rightarrow \tau(B)$ is continuous linear map between two Banach spaces. Thus it is a homeomorphism.

□

Corollary 8. *Let X be a Banach space and Y is a closed linear subspace of X .*

- 1) *X is reflexive if and only if the unit ball in X is weakly compact.*
- 2) *If X is reflexive and $x_0 \in X \setminus Y$, then there is a point $y_0 \in Y$ s.t. $\|x_0 - y_0\| = \inf\{y - x_0 : y \in Y\}$.*

Proof. For 1), by above theorem, the unit ball in X^{**} is wk^* -compact. Since $X = X^{**}$, the unit ball in X is wk -compact. The converse is similar as above proof.

For 2), Y is wk -compact by above corollary. Since by the Hahn-Banach Theorem, for each x , we can find a $f \in X^*$ with $\|f\| = 1$ and $f(x) = x$. The map $x \mapsto \|x - x_0\|$ is weakly lower semicontinuous. Therefore, there is a $y_0 \in Y$ s.t. $\|x_0 - y_0\| = \inf\{y - x_0 : y \in Y\}$. □

2.4.4 Adjoint Operators

Definition 1. Let $T: X \rightarrow Y$ be a bounded linear map between two Banach spaces. Then the adjoint of operator of T is defined as

$$\begin{aligned} T^*: Y^* &\longrightarrow X^* \\ f &\longmapsto f \circ T \end{aligned}$$

There are some easy properties of the adjoint operator, which can be obtained by the definition.

Proposition 2. *Let X, Y and Z be three Banach spaces and $T \in \mathcal{B}(X, Y), S \in \mathcal{B}(Y, Z)$. Then, we have*

- 1) $T^{**}|_X = T$
- 2) the map $T \rightarrow T^*$ is an isometric isomorphism.
- 3) $(ST)^* = T^*S^*$.
- 4) $\ker T^* = (\text{ran } T)^\perp$ and $\ker T = (\text{ran } T^*)^\perp$.

Then for the 4) in above theorem, we can find dual consequences. Firstly, there is a useful lemma.

Lemma 3. Let $T: X \rightarrow Y$ be a bounded linear map between two Banach spaces. Then $\text{ran } T$ is closed if and only if there exists a constant $C > 0$, for any $y \in \text{ran } T$ there is a point $x \in X$ s.t. $\|y\| \geq C \|x\|$.

Proof. By the Open Mapping Theorem, for the unit open ball $B \subset X$ and some $\delta > 0$, s.t.

$$\{y \in \text{ran } T: \|y\| < \delta\} \subset T(B)$$

Thus for any nonzero $y \in \text{ran } T$,

$$\exists z \in B \text{ s.t. } Tz = \frac{\delta}{2\|y\|}y \Rightarrow \|y\| \geq \frac{2}{\delta}\|x\| \text{ where } x = \frac{2\|y\|}{\delta}z \text{ with } Tx = y$$

Conversely, if $y \in \overline{\text{ran } T}$, there is $\{y_n\} \subset \text{ran } T$ with $y_n \rightarrow y$. By the assumption, there is a constant $C > 0$ and a sequence $\{x_n\} \subset X$ s.t. $\|x_n - x_m\| \leq C \|Tx_n - Tx_m\|$. Since $\{y_n\}$ is Cauchy, $\|x_n - x_m\|$ is Cauchy. Thus there exists a $x \in X$ s.t. $x_n \rightarrow x$ and $Tx = y$. \square

Remark. There is a more special case for this lemma. $\ker T = \{0\}$ and $\text{ran } T$ is closed if and only if there exists a constant $C > 0$ s.t. $\|Tx\| \geq C \|x\|$ for any $x \in X$.

Theorem 4. Let $T: X \rightarrow Y$ be a bounded linear map between two Banach spaces.

$$1) \overline{\text{ran } T} = (\ker T^*)^\perp.$$

$$2) \text{ If } \text{ran } T \text{ is closed, then } \text{ran } T^* \text{ is closed and } \text{ran } T^* = (\ker T)^\perp.$$

Proof. For 1), let $y \in \overline{\text{ran } T}$ and $y_n \in \text{ran } T$ with $y_n \rightarrow y$. If $g \in \ker T^*$, then

$$g(y_n) = g(Tx_n) = T^*g(x_n) = 0 \rightarrow g(y) \text{ i.e. } y \in (\ker T^*)^\perp$$

Conversely, if $y_0 \notin \overline{\text{ran } T}$, then by the Hahn-Banach Theorem, there is a $g \in Y^*$ s.t.

$$\text{ran } T \subset \ker g \text{ i.e. } g \in \ker T^* \text{ \& } g(y_0) \neq 0 \Rightarrow y_0 \notin (\ker T^*)^\perp$$

For 2), if $f \in (\ker T)^\perp$ i.e. $\ker T \subset \ker f$, then there is a $\tilde{g} \in (\text{ran } T)^*$ s.t. $\tilde{g}(Tx) = f(x)$ since $\text{ran } T \cong X/\ker T$. By above lemma, there exists a constant $C > 0$ s.t. for any $y \in \text{ran } T$ there is a $x \in X$ with $Tx = y$ and $\|x\| \leq C \|y\|$. Therefore,

$$|\tilde{g}(y)| = |f(x)| \leq C \|f\| \|x\|$$

Thus \tilde{g} can extend to g defining on Y^* s.t.

$$T^*g(x) = g(Tx) = f(x) \text{ i.e. } T^*g = f$$

Therefore, $(\ker T)^\perp \subset \text{ran } T^*$.

Conversely, if $f \in \text{ran } T^*$, i.e. there is a $g \in Y^*$ with $T^*g = f$, then for any $x \in \ker T$, $f(x) = T^*g(x) = g(Tx) = 0$. Therefore, $f \in (\ker T)^\perp$. \square

Corollary 5. *Let $T: X \rightarrow Y$ be a bounded linear map between two Banach spaces. Then T is invertible if and only if T^* is invertible. In this case, $(T^*)^{-1} = (T^{-1})^*$.*

2.5 Hilbert Spaces

A Hilbert space is a special Banach space, which is endowed with an inner product. And the structure of inner product can provide better properties of Hilbert spaces.

2.5.1 Projection Theorem and Riesz Theorem

Definition 1. On a vector space \mathcal{H} over \mathbb{C} , an inner product is a map $\langle \cdot, \cdot \rangle: \mathcal{H} \times \mathcal{H} \rightarrow \mathbb{K}$ satisfies that for any $\alpha, \beta \in \mathbb{K}$ and $x, y, z \in \mathcal{H}$,

- 1) $\langle \alpha x + \beta y, z \rangle = \alpha \langle x, z \rangle + \beta \langle y, z \rangle$
- 2) $\langle x, y \rangle = \overline{\langle y, x \rangle}$
- 3) $\langle x, x \rangle \geq 0$ and $\langle x, x \rangle = 0 \Leftrightarrow x = 0$

Remark. The vector space \mathcal{H} with the inner product $\langle \cdot, \cdot \rangle$ is called an inner product space. And the inner product can induce the CBS-Inequality,

$$|\langle x, y \rangle|^2 \leq \langle x, x \rangle \langle y, y \rangle$$

by using the positivity of the inner product, $\langle x - \alpha y, x - \alpha y \rangle \geq 0$.

Then on the \mathcal{H} , a nature norm can be induced by the inner product, defined as $\|x\|^2 = \langle x, x \rangle$ for any $x \in \mathcal{H}$. This definition is valid since the linearity, positivity and CBS-Inequality of the inner product. And there are two forms of the norm.

Proposition 2. *Let $(\mathcal{H}, \langle \cdot, \cdot \rangle)$ be an inner product space, and $\|\cdot\|$ be the coincided with this inner product. Then, we have the following identities.*

Polarization identity:

$$\langle x, y \rangle = \frac{1}{4}(\|x + y\|^2 - \|x - y\|^2 + i\|x - iy\|^2 - i\|x + iy\|^2)$$

Parallelogram law:

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

Remark. The first identity can be used to construct an inner product, but not all norm can do this successfully. So, the second identity can be used to check whether a norm can construct an inner product on a normed space.

Definition 3. A Hilbert space is an inner product space, which is complete with respect to the induced norm.

Therefore, a Hilbert space is indeed a Banach space, whose norm can construct an inner product. Now, we can find the extra properties provided by the inner product. The following theorem is the most important property the Hilbert space has.

Theorem 4. If \mathcal{H} is a Hilbert space and K is a closed convex nonempty subset of \mathcal{H} and $x \in \mathcal{H} \setminus K$, then there exists a unique $k_0 \in K$ s.t.

$$\|x - k_0\| = \inf \|x - k\| : k \in K$$

Proof. Let $d = \inf \|x - k\| : k \in K$, then there is a sequence $\{k_n\} \subset K$ s.t. $\|x - k_n\| \rightarrow d$. By the parallelogram law,

$$4 \left\| x - \frac{1}{2}(k_m + k_n) \right\|^2 + \|k_m - k_n\|^2 = 2(\|x - k_m\|^2 + \|x - k_n\|^2) \rightarrow 4d^2$$

Since $\frac{1}{2}(k_m + k_n) \in K$,

$$4 \left\| x - \frac{1}{2}(k_m + k_n) \right\|^2 \geq 4d^2$$

Thus $\|k_m - k_n\| \rightarrow 0$. By the facts that \mathcal{H} is complete and K is closed, there exists $k_0 \in K$ s.t. $k_n \rightarrow k_0$. Then $\|x - k_0\| = \lim_{n \rightarrow \infty} \|x - k_n\| = d$. If there is another $k_1 \in K$, $\|x - k_0\| = \|x - k_1\|$.

$$d \leq \left\| x - \frac{1}{2}(k_0 + k_1) \right\| \leq \frac{1}{2}(\|x - k_0\| + \|x - k_1\|) = 2d$$

Then

$$\left\| x - \frac{1}{2}(k_0 + k_1) \right\| = \frac{1}{2}(\|x - k_0\| + \|x - k_1\|)$$

Thus by the parallelogram law, $k_0 = k_1$. □

Remark. The main method of above proof is the parallelogram law, which is the most essential property of the norm induced by the inner product.

For $x, y \in \mathcal{H}$, if $\langle x, y \rangle = 0$, we say x is orthogonal to y . Similarly, if W is a subset of \mathcal{H} , we define $W^\perp = \{y \in \mathcal{H} : \langle x, y \rangle = 0 \forall x \in W\}$. Then we can use above theorem to obtain an important structure of a Hilbert space.

Theorem 5 (Projection Theorem). Let \mathcal{M} be a closed subspace of a Hilbert space H . Then $\mathcal{H} = \mathcal{M} \oplus \mathcal{M}^\perp$.

Proof. Clearly, $\mathcal{M} \cap \mathcal{M}^\perp = \{0\}$. We just need to prove that for any $x \in \mathcal{H}$, there exist a unique $x_1 \in \mathcal{M}$ and $x_2 \in \mathcal{M}^\perp$, s.t. $x = x_1 + x_2$. If $x \in \mathcal{M}$, let $x_1 = x$ and $x_2 = 0$.

Assume $x \notin \mathcal{M}$, then by above theorem, we can find a unique $m_0 \in \mathcal{M}$ s.t. $\|x - m_0\| = \inf \{\|x - m\| : m \in \mathcal{M}\}$. Let $x_1 = m_0$. Then $x = x_1 + (x - m_0)$. Therefore, it is sufficient to prove that $x - m_0 \in \mathcal{M}^\perp$.

For any $m \in \mathcal{M}$ and any $\lambda \in \mathbb{C}$, since $m_0 + \lambda m \in \mathcal{M}$, we have

$$\begin{aligned} \|x - m_0\|^2 &\leq \|x - m_0 - \lambda m\|^2 \\ &= \|x - m_0\|^2 - 2 \operatorname{Re} \lambda \langle m, x - m_0 \rangle + |\lambda|^2 \|m\|^2 \end{aligned}$$

Thus $-2 \operatorname{Re} \lambda \langle m, x - m_0 \rangle + |\lambda|^2 \|m\|^2 \geq 0$. Then taking the $\lambda = \varepsilon > 0$ and let $\varepsilon \rightarrow 0$, therefore

$$\operatorname{Re} \langle m, x - m_0 \rangle \leq 0$$

Similarly, taking $\lambda = -i\varepsilon$ and let $\varepsilon \rightarrow 0$,

$$\operatorname{Im} \langle m, x - m_0 \rangle \leq 0$$

Also, these are true for $-m$. Therefore, $\langle m, x - m_0 \rangle = 0$ for any $m \in \mathcal{M}$, i.e. $x - m_0 \in \mathcal{M}^\perp$. \square

By the Projection Theorem, we can directly obtained the following corollary.

Corollary 6. *Let \mathcal{H} be a Hilbert space.*

- 1) *If \mathcal{M} is a closed subspace of \mathcal{H} , then $(\mathcal{M}^\perp)^\perp = \mathcal{M}$.*
- 2) *If $A \subset \mathcal{H}$ is a subset, then $(A^\perp)^\perp = \overline{\operatorname{span} A}$.*

Using the Projection Theorem, the dual space of a Hilbert space can be more explicit.

Theorem 7 (Riesz Theorem). *Let \mathcal{H} be a Hilbert space. Then the map σ is defined as*

$$\begin{aligned} \sigma: \mathcal{H} &\longrightarrow \mathcal{H}^* \\ x &\longmapsto L_x \end{aligned}$$

where $L_x(y) = \langle y, x \rangle$. Then σ is an isometric antilinear bijection, i.e. $\mathcal{H} \cong \mathcal{H}^*$.

Proof. Clearly, L_x is antilinear. And by the CBS-Inequality, $\|L_x\| = \|x\|$, thus σ is definitely an isometry. Then we just need to prove σ is surjective. If $L \in \mathcal{H}^*$ is nonzero, then there is a $x_0 \in (\ker L)^\perp$. Thus we can assume $L(x_0) = 1$. If $y \in \mathcal{H}$, then $y - L(y)x_0 \in \ker L$. Therefore,

$$0 = \langle y - L(y)x_0, x_0 \rangle = \langle y, x_0 \rangle - L(y) \|x_0\|^2$$

Then let $x = \|x_0\|^{-2} x_0$, $L_x = L$. \square

Corollary 8. *A Hilbert space \mathcal{H} is reflexive. Thus, \mathcal{H} is weakly complete and a subset in \mathcal{H} is weakly compact if and only if it is bounded and weakly closed.*

And the Riesz Theorem can extend to bounded sesquilinear forms.

Definition 9. Let \mathcal{H}_1 and \mathcal{H}_2 be two Hilbert spaces. The map

$$f: \mathcal{H}_1 \times \mathcal{H}_2 \rightarrow \mathbb{C}$$

is called a sesquilinear form, if it satisfies

- 1) $f(\alpha x + \beta y, z) = \alpha f(x, z) + \beta f(y, z)$
- 2) $f(x, \alpha y + \beta z) = \bar{\alpha} f(x, y) + \bar{\beta} f(x, z)$

Remark. If f is continuous, then we know $|f(x, y)| \leq C \|x\| \|y\|$ for some $C < \infty$. Also, the converse is true. Then we can define the norm of f as

$$\|f\| = \sup_{x \in \mathcal{H}_1, y \in \mathcal{H}_2} \frac{|f(x, y)|}{\|x\| \|y\|}$$

Theorem 10. Let $f: \mathcal{H}_1 \times \mathcal{H}_2 \rightarrow \mathbb{C}$ be a bounded sesquilinear form of two Hilbert spaces \mathcal{H}_1 and \mathcal{H}_2 . Then there is a unique bounded linear map with $\|S\| = \|f\|$

$$S: \mathcal{H}_1 \rightarrow \mathcal{H}_2 \text{ s.t. } f(x, y) = \langle Sx, y \rangle \quad \forall x \in \mathcal{H}_1 \quad \forall y \in \mathcal{H}_2$$

2.5.2 Adjoint Operators

Let $T: \mathcal{H} \rightarrow \mathcal{H}$ be a bounded linear operator on a Hilbert space \mathcal{H} . Since \mathcal{H} is also a Banach space, we can define the adjoint operator of T as

$$\begin{aligned} T^*: \mathcal{H}^* &\longrightarrow \mathcal{H}^* \\ f &\longmapsto T^* f \end{aligned}$$

By the Riesz Theorem, each $f \in \mathcal{H}^*$ can be expressed as $f = \langle \cdot, x \rangle$ for some $x \in \mathcal{H}$. Then for any $y \in \mathcal{H}$ we have

$$T^*(\langle \cdot, x \rangle)(y) = (\langle \cdot, x \rangle)(Ty) = \langle Ty, x \rangle$$

Clearly, $f(y) = \langle Ty, x \rangle \in \mathcal{H}^*$, thus there is a unique $\tilde{x} \in \mathcal{H}$ s.t. $f(y) = \langle y, \tilde{x} \rangle$. Then, we have

$$T^*(\langle \cdot, x \rangle)(y) = \langle Ty, x \rangle = \langle y, \tilde{x} \rangle = (\langle \cdot, \tilde{x} \rangle)(y)$$

Thus, we have

$$T^*(\langle \cdot, x \rangle) = \langle \cdot, \tilde{x} \rangle$$

Since $\mathcal{H} \cong \mathcal{H}^*$, in fact $T^*: \mathcal{H} \rightarrow \mathcal{H}$, and by above mention, $T^*(x) = \tilde{x}$, where \tilde{x} is determined by the equation

$$\langle y, \tilde{x} \rangle = \langle Ty, x \rangle, \quad \forall y \in \mathcal{H}$$

Therefore, we can give the definition of the adjoint operator on a Hilbert space.

Definition 1. Let $T: \mathcal{H} \rightarrow \mathcal{H}$ be a bounded linear operator on a Hilbert space \mathcal{H} . Then the adjoint operator $T^*: \mathcal{H} \rightarrow \mathcal{H}$ is defined as

$$\langle y, T^*x \rangle = \langle Ty, x \rangle \quad \forall y \in \mathcal{H}$$

Remark. Equivalently, $\langle T^*x, y \rangle = \langle x, Ty \rangle \quad \forall y \in \mathcal{H}$. By above mention, this definition is well-defined and nature.

Since this definition is induced by the definition of a Banach space, thus all results in the subsection 2.4.4 also hold for the adjoint operators on a Hilbert space, except the $T \rightarrow T^*$ is linear. In fact, on a Hilbert space, this map $T \rightarrow T^*$ is antilinear, i.e. $((\alpha + \beta)T)^* = \bar{\alpha}T^*T^* + \bar{\beta}T^*$. Moreover, Since any Hilbert space is reflexive, the first result in **Proposition 2** in the subsection 1.4.4 can be $T^{**} = T$.

And by the **Corollary 6** of the Projection Theorem and the last result in **Proposition 2** in the subsection 1.4.4, we have the following proposition.

Proposition 2. Let $T: \mathcal{H} \rightarrow \mathcal{H}$ be a bounded linear operator on a Hilbert space \mathcal{H} . Then, we have

$$\begin{aligned} \ker T^* &= (\text{ran } T)^\perp, & \ker T &= (\text{ran } T^*)^\perp \\ \overline{\text{ran } T} &= (\ker T^*)^\perp, & \overline{\text{ran } T^*} &= (\ker T)^\perp \end{aligned}$$

And there is a result, which provide another important information of operators on a Hilbert space and make it more interesting than the operators defined on a Banach space.

Theorem 3. Let $T \in \mathcal{B}(\mathcal{H})$ for a Hilbert space. Then

$$\|T\| = \|T^*\| = \|T^*T\|^{\frac{1}{2}}$$

Proof. Firstly, for any $h \in \mathcal{H}$ with $\|h\| \leq 1$

$$\|Ah\|^2 = \langle Ah, Ah \rangle = \langle A^*Ah, h \rangle \leq \|A^*Ah\| \|h\| \leq \|A^*A\| \|h\| \leq \|A^*\| \|A\|$$

Therefore, $\|A\|^2 \leq \|A^*A\| \leq \|A^*\| \|A\|$. and $\|A\| \leq \|A^*\|$. But $(A^*)^* = A$, thus this identity can hold. \square

2.5.3 Orthonormal Sets and Schauder Basis

For a Hilbert space \mathcal{H} , subset $W \subset \mathcal{H}$ is called orthonormal if elements in W are pairwise orthogonal and each element has norm 1. Then we can find a maximal orthonormal set, which is called a Schauder basis. Since the algebraic basis, Hamel basis, of a Hilbert space may be uncountable, we firstly need to define the uncountable summation.

Definition 1. Let $\{h_i : i \in I\}$ be family of elements in \mathcal{H} and \mathcal{F} be a collection of all finite subsets of I . \mathcal{F} can be endowed with the order by \subset , then \mathcal{F} is a directed set. Define the net, where $F \in \mathcal{F}$

$$h_F = \sum \{h_i : i \in F\}$$

Therefore, we can define sum of $\{h_i : i \in I\}$ as

$$\sum \{h_i : i \in I\} = \lim h_F$$

where the limit is with respect to the norm topology on H .

For finite orthonormal set, the results can be obtained by the algebraic structure. And if the orthonormal set is infinite, we have similar results.

Theorem 2 (Bessel's Inequality). *If $\{e_n\}_{n=1}^{\infty}$ is an orthonormal set of a Hilbert space and $h \in \mathcal{H}$, then*

$$\sum_{n=1}^{\infty} |\langle h, e_n \rangle|^2 \leq \|h\|^2$$

Proof. Let $h_n = h - \sum_{k=1}^n \langle h, e_k \rangle e_k$. Then h_n is clearly orthogonal to $\{e_k\}_{k=1}^n$.

$$\begin{aligned} \|h\|^2 &= \|h_n\|^2 + \left\| \sum_{k=1}^n \langle h, e_k \rangle e_k \right\|^2 \\ &= \|h_n\|^2 + \sum_{k=1}^n |\langle h, e_k \rangle|^2 \\ &\geq \sum_{k=1}^n |\langle h, e_k \rangle|^2 \end{aligned}$$

□

Corollary 3. *Let \mathcal{H} be a Hilbert space and \mathcal{E} be an orthonormal subset.*

- 1) *For any $h \in \mathcal{H}$, $|\langle h, e \rangle| \neq 0$ for at most a countable many $e \in \mathcal{E}$.*
- 2) $\sum_{e \in \mathcal{E}} |\langle h, e \rangle|^2 \leq \|h\|^2$.

3) $\sum \{\langle h, e \rangle e : e \in \mathcal{E}\}$ is converges.

Proof. For 1), let $\mathcal{E}_n = \{e \in \mathcal{E} : |\langle h, e \rangle| \geq \frac{1}{n}\}$, then \mathcal{E}_n has to be finite. But

$$\bigcup_{n=1}^{\infty} \mathcal{E}_n = \{e \in \mathcal{E} : |\langle h, e \rangle| \neq 0\}$$

For 2), it is clearly by the Bessel's Inequality and above corollary.

For 3), it is because the net $h_F = \sum \{h_i : i \in F\}$, where $F \subset \mathcal{E}$ is finite, is Cauchy by using the fact that for any $\varepsilon > 0$ there is a $N \in \mathbb{N}$ s.t.

$$\sum_{n=N}^{\infty} |\langle h, e_n \rangle| < \varepsilon$$

□

Definition 4 (Schauder basis). A orthonormal set \mathcal{E} of a Hilbert space \mathcal{H} is called a Schauder basis if for any $h \in \mathcal{H}$,

$$h = \sum \{\langle h, e \rangle e : e \in \mathcal{E}\}$$

Remark. By above corollary, this definition is equivalent to that for any $h \in \mathcal{H}$, there exists $\{e_n\}_{n=1}^{\infty} \subset \mathcal{E}$, s.t.

$$h = \sum_{n=1}^{\infty} \langle h, e_n \rangle e_n$$

Moreover, by the Zorn's Lemma, every Hilbert space has a Schauder basis.

By the definition, and above theorem and corollaries, we can find the properties of the Schauder basis.

Theorem 5. Let \mathcal{H} be a Hilbert space and \mathcal{E} be an orthonormal subset. Then the following statements are equivalent.

- 1) \mathcal{E} is a Schauder basis.
- 2) $h \in \mathcal{H}$ with $h \perp \mathcal{E}$, then $h = 0$.
- 3) $g, h \in \mathcal{H}$, then

$$\langle g, h \rangle = \sum \{\langle g, e \rangle \langle e, h \rangle : e \in \mathcal{E}\}$$

- 4) (Parseval's Identity) $h \in \mathcal{H}$, then

$$\|h\|^2 = \sum \{|\langle h, e \rangle|^2 : e \in \mathcal{E}\}$$

Proposition 6. Let \mathcal{H} be a Hilbert space and \mathcal{E} be an orthonormal subset.

- 1) \mathcal{H} is separable if and only if \mathcal{E} is countable.
- 2) Any separable Hilbert space is isometrically isomorphic to $l^2(\mathbb{C})$.

2.5.4 Unitaries and Projections

If we can find a bijection, which is continuous and the inverse of which is also continuous, between two topological spaces, these two topological spaces are regarded as same. Also, if there is a bijection, which is linear between two linear spaces, then we say these two linear spaces are same. Now, we also want to use a special map to classify Hilbert spaces. We know the Hilbert space has two structures, linear structure and the inner product structure. So following definition is nature.

Definition 1. Let \mathcal{H} and \mathcal{K} be two Hilbert spaces. If there is a linear isomorphism U from \mathcal{H} to \mathcal{K} , s.t $\langle Ux, Uy \rangle = \langle x, y \rangle$, then U is called a unitary.

Remark. In fact, if there exists a unitary U between \mathcal{H} and \mathcal{K} , U preserves not only the linear structure, but also the inner product structure of Hilbert spaces. Therefore, these two Hilbert spaces are regarded as same. Moreover, if the U is not assumed bijection, then U is called an isometry.

As above definition, the following proposition is clearly.

Proposition 2. If $U \in \mathcal{B}(\mathcal{H})$, then the following statements are equivalent.

- 1) U is a unitary.
- 2) $U^*U = UU^* = I$.
- 3) U is a normal isometry.

Remark. $A \in \mathcal{B}(\mathcal{H})$ is called normal, if $A^*A = AA^*$.

Now, the another important operator is projection. Like its name showing, the projection can project the elements in a Hilbert space to a corresponding subspace, and because of the Projection Theorem, we can make the definition be well-defined.

Definition 3. Let \mathcal{M} be a any closed subspace of \mathcal{H} . Then by the Projection Theorem, $\mathcal{H} = \mathcal{M} \oplus \mathcal{M}^\perp$, i.e. for any $x \in \mathcal{H}$, there exist un $x_{||} \in \mathcal{M}$ and $x_\perp \in \mathcal{M}^\perp$ s.t. $x = x_{||} + x_\perp$. Therefor we can define a map $P: \mathcal{H} \rightarrow \mathcal{M}$ as $Px = x_{||}$. P is called a projection.

Remark. P is a linear operator can be easily proved.

There is an equivalent definition of the projection. Firstly, $A \in \mathcal{B}(\mathcal{H})$ is called self-adjoint if $A^* = A$.

Theorem 4. Let \mathcal{H} be a Hilbert space. $P \in \mathcal{B}(\mathcal{H})$ is a projection if and only if P is self-adjoint and $P^2 = P$.

Proof. Let $\mathcal{H} = \mathcal{M} \oplus \mathcal{M}^\perp$. If P is a projection to \mathcal{M} , $P^2 = P$ is clearly true.

$$\begin{aligned}\langle P(x = x_{||} + x_\perp), y = y_{||} + y_\perp \rangle &= \langle x_{||}, y_{||} \rangle \\ \langle P^*(x = x_{||} + x_\perp), y = y_{||} + y_\perp \rangle &= \langle x = x_{||} + x_\perp, P(y = y_{||} + y_\perp) \rangle \\ &= \langle x_{||}, y_{||} \rangle\end{aligned}$$

Therefore, P is self-adjoint.

Conversely, $\mathcal{M} = P(\mathcal{H})$ is clearly a subspace. And for any $y \in \mathcal{M}$, let $x = y$, we have $Px = y$. Thus by the **Lemma 3** in the subsection 1.4.4, \mathcal{M} is closed. Since P is self-adjoint,

$$\langle x - Px, Px \rangle = \langle Px - P^2x, x \rangle = 0$$

That means $x - Px \in \mathcal{M}^\perp$. Therefore, $x = Px + (x - Px)$ is the unique decomposition of x with respect to \mathcal{M} . P is a projection. \square

Remark. By the fact that $P^2 = P$ and the **Theorem 3** in the subsection 1.5.2, $\|P\| = 1$. Therefore, P is bounded.

The one of the importance of projections is that any closed subspace of a Hilbert space can be regarded as a operator on the Hilbert space. And then this can provide us an effective method to research the invariant space.

Definition 5. For $T \in \mathcal{B}(\mathcal{H})$ and a closed subspace \mathcal{M} of \mathcal{H} , if $T(\mathcal{M}) \subset \mathcal{M}$, \mathcal{M} is called an invariant space for T . If $T(\mathcal{M}) \subset \mathcal{M}$ and $T(\mathcal{M}^\perp) \subset \mathcal{M}^\perp$, then \mathcal{M} is called a reducing space for T .

Proposition 6. Let \mathcal{M} be a closed subspace of a Hilbert space \mathcal{H} and P be the projection to \mathcal{M} and $T \in \mathcal{B}(\mathcal{H})$. T acts on $\mathcal{H} = \mathcal{M} \oplus \mathcal{M}^\perp$ can be expressed as

$$T = \begin{pmatrix} W & X \\ Y & Z \end{pmatrix}$$

Then the following statements are equivalent.

- 1) \mathcal{M} is invariant for T .
- 2) $PTP = TP$.
- 3) $Y = 0$.

Proof. 1) \Rightarrow 2) is clearly.

2) \Rightarrow 3): The P can be expressed as

$$P = \begin{pmatrix} I & 0 \\ 0 & 0 \end{pmatrix}$$

Then by $PTP = TP$, we can find $Y = 0$.

3) \Rightarrow 1) is also clearly. \square

Proposition 7. *Let \mathcal{M}, \mathcal{H} and T, P be defined as above proposition. Then the following statements are equivalent.*

1) \mathcal{M} reduces \mathcal{H} .

2) $PT = TP$.

3) X and Y are 0.

4) \mathcal{M} is invariant for T and T^* .

Proof. 1) \Rightarrow 2): By above proposition, $PTP = TP$ and $(I - P)T(I - P) = T(I - P)$, thus we have $PT = TP$.

2) \Rightarrow 3): Expressing P as above proposition, then by $PT = TP$, we find X and Y are 0.

3) \Rightarrow 4): Since X and Y are 0, then

$$T^* = \begin{pmatrix} W^* & 0 \\ 0 & Z^* \end{pmatrix}$$

Clearly, \mathcal{M} is invariant for T and T^* .

4) \Rightarrow 1): If $h \in \mathcal{M}^\perp$ and $g \in \mathcal{M}$, $\langle g, Ah \rangle = \langle A^*g, h \rangle = 0$ since $A^*g \in \mathcal{M}$. Therefore, $T(\mathcal{M}^\perp) \subset \mathcal{M}^\perp$. \square

Chapter 3

Banach Algebras and C^* -Algebras

3.1 Banach Algebras

Let \mathcal{H} be a Hilbert space. Then $\mathcal{B}(\mathcal{H})$ is indeed a Banach space. But we have more structure on it. Any two element $S, T \in \mathcal{B}(\mathcal{H})$ can do multiplication, defined as $ST = S \circ T$, then $ST \in \mathcal{B}(\mathcal{H})$ by the definition of norm, and moreover $\|ST\| \leq \|S\| \|T\|$. Therefore, $\mathcal{B}(\mathcal{H})$ is an algebra with the extra property of the multiplication, called a Banach algebra.

3.1.1 Elementary Properties

Definition 1. A Banach algebra \mathfrak{A} is an algebra over \mathbb{C} with a norm $\|\cdot\|$ relative to which \mathfrak{A} is a Banach space and s.t. for all $a, b \in \mathfrak{A}$,

$$\|ab\| \leq \|a\| \|b\|.$$

Remark. The extra condition guarantees the multiplication is norm continuous. In fact, \mathfrak{A} with the multiplication and the norm is a topological semigroup.

If \mathfrak{A} has an identity 1, it assumes $\|1\| = 1$. But if \mathfrak{A} does not have an identity, we can add an identity to it.

Proposition 2. If \mathfrak{A} is a Banach algebra without the identity. The let $\tilde{\mathfrak{A}} = \mathfrak{A} \oplus \mathbb{C}$ be an induced vector space. Then we define a norm on $\tilde{\mathfrak{A}}$ as

$$\|(a, \lambda)\| = \|a\| + |\lambda|$$

and define a multiplication on $\tilde{\mathfrak{A}}$ as

$$(a, \alpha)(b, \beta) = (ab + \alpha b + \beta a, \alpha\beta)$$

Then $\tilde{\mathfrak{A}}$ is a Banach algebra with the identity $(0, 1)$.

Then \mathfrak{A} can be isometrically imbedded into $\tilde{\mathfrak{A}}$. Therefore, we can always assume \mathfrak{A} has an identity. For a unit algebra, the invertibility of a element is important.

Theorem 3. *Let \mathfrak{A} be a Banach algebra and $a \in \mathfrak{A}$. If $\|a - 1\| < 1$, then a is invertible.*

Proof. For a nonzero real number $\lambda \in \mathbb{R}$ with $|\lambda| < 1$, we know

$$(1 - \lambda)^{-1} = \sum_{n=0}^{\infty} \lambda^n$$

Similar, for $a \in \mathfrak{A}$ with $\|a - 1\| < 1$, set

$$b = \sum_{n=0}^{\infty} (1 - a)^n$$

Firstly, since $\|(1 - a)^n\| \leq \|(1 - a)\|^n$, b is well-defined. Then we can prove $b = (1 - (1 - a))^{-1} = a^{-1}$. \square

Remark. This result is important. It says a small perturbation of an invertible element is also invertible. It is because that the multiplication is norm continuous. And by the continuity of multiplication, this result can be true at any point other than 1.

Corollary 4. *Let \mathfrak{A} be a Banach algebra and*

$$\begin{aligned} G_l &= \{ a \in \mathfrak{A} : a \text{ is left invertible} \} \\ G_r &= \{ a \in \mathfrak{A} : a \text{ is right invertible} \} \\ G &= G_l \cap G_r = \{ a \in \mathfrak{A} : a \text{ is invertible} \} \end{aligned}$$

Then G_l and G_r and G are open. Moreover, the map $a \rightarrow a^{-1}$ from G to G is continuous.

Proof. G_l and G_r and G are open by above theorem.

We just need to check this map is continuous at 1 because of the continuity of multiplication. For $\{a_n\} \subset G$ with $a_n \rightarrow 1$, thus $\|1 - a_n\| < \delta < 1$. Since

$$a_n^{-1} = (1 - (1 - a_n))^{-1} = \sum_{k=0}^{\infty} (1 - a_n)^k = 1 + \sum_{k=1}^{\infty} (1 - a_n)^k$$

Therefore, we have

$$\begin{aligned} \|1 - a_n^{-1}\| &= \left\| \sum_{k=1}^{\infty} (1 - a_n)^k \right\| \\ &\leq \sum_{k=1}^{\infty} \|1 - a_n\|^k \\ &< \frac{\delta}{1 + \delta} < \delta = \|1 - a_n\| \end{aligned}$$

i.e. $\lim a_n = 1$. □

Corollary 5. *Let \mathfrak{A} be a Banach algebra.*

- 1) *The closure of a proper ideal is a proper ideal.*
- 2) *A maximal ideal is closed.*
- 3) *every ideal contained in a maximal ideal.*

If \mathfrak{B} is a closed ideal of a Banach algebra \mathfrak{A} , then the quotient algebra $\mathfrak{A}/\mathfrak{B}$ with the induced norm is also a Banach algebra since

$$\|(a + \mathfrak{B})(b + \mathfrak{B})\| = \|ab + \mathfrak{B}\| \leq \|(a + b_1)(b + b_2)\| \leq \|(a + b_1)\| \|(b + b_2)\|$$

for any $b_1, b_2 \in \mathfrak{B}$.

3.1.2 Spectrum

Definition 1. Let \mathfrak{A} be a Banach algebra and $a \in \mathfrak{A}$. The spectrum of a , denoted by $\sigma(a)$ defined as

$$\sigma(a) = \{ \lambda \in \mathbb{C} : a - \lambda \text{ is invertible} \}$$

And the resolvents of a , $\rho(a) = \mathbb{C} \setminus \sigma(a)$.

Moreover, we can define the spectral radius of a as

$$r(a) = \sup \{ |\lambda| : \lambda \in \sigma(a) \}$$

Firstly, there are some elementary properties of the spectrum.

Theorem 2. *Let \mathfrak{A} be a Banach algebra and $a \in \mathfrak{A}$.*

- 1) *If $|\lambda| > \|a\|$, then $\lambda \notin \sigma(a)$.*
- 2) *$\sigma(a)$ is a compact subset of \mathbb{C} .*
- 3) *the map $\lambda \mapsto (a - \lambda)^{-1}$ from $\rho(a)$ to \mathfrak{A} is analytic and $\sigma(a)$ is nonempty.*
- 4) *$r(a) = \lim_{n \rightarrow \infty} \|a^n\|^{\frac{1}{n}}$.*

Proof. 1) holds by above theorem.

For 2), since $\lambda \rightarrow a - \lambda$ is continuous from \mathbb{C} to \mathfrak{A} and G is open, $\rho(a)$ is open i.e. $\sigma(a) = \mathbb{C} \setminus \rho(a)$ is closed. Then by 1), $\sigma(a)$ is compact.

For 3), by the identity $a^{-1} - b^{-1} = a^{-1}(b - a)b^{-1}$ and the continuity of $a \rightarrow a^{-1}$, we can compute the derivative of $F(\lambda) = (a - \lambda)^{-1}$,

$$F'(\lambda) = (a - \lambda)^{-2}$$

And clearly, $F'(\lambda)$ is continuous. Thus it is analytic and it vanishes at ∞ . By the Liouville's Theorem, if $\rho(a) = \mathbb{C}$, F is constant. Therefore, $\rho(a) \neq \mathbb{C}$ i.e. $\sigma(a) \neq \emptyset$.

For 4), let $U = \{\lambda \in \mathbb{C} : \lambda = 0 \text{ or } \lambda^{-1} \in \rho(a)\}$ and

$$f(\lambda) = \begin{cases} (\lambda^{-1} - a)^{-1} & x \neq 0, \\ 0, & x = 0 \end{cases}$$

Then f is analytic on U , i.e. $f(\lambda) = \lambda \sum_{n=0}^{\infty} \lambda^n a^n$ is well-defined. Therefore, the convergent radius $R = r(a)^{-1}$

$$R^{-1} = \limsup_{n \rightarrow \infty} \|a^n\|^{\frac{1}{n}} = r(a)$$

Conversely, by the identity $(a^n - \lambda^n) = (a - \lambda)(a^{n-1} + \lambda a^{n-2} + \lambda^2 a^{n-3} + \dots + \lambda^{n-1})$. Then, if $(a^n - \lambda^n)$ is invertible, then $(a - \lambda)$ is invertible, i.e. $\sigma(a) \subset \sigma(a^n)$. Thus $|\lambda|^n \leq \|a^n\|$ for any $\lambda \in \sigma a$. $r(a) = \liminf_{n \rightarrow \infty} \|a^n\|^{\frac{1}{n}}$. Therefore, $r(a) = \lim_{n \rightarrow \infty} \|a^n\|^{\frac{1}{n}}$. \square

If $\mathfrak{B} \subset \mathfrak{A}$ is a subalgebra with the same identity of a Banach algebra \mathfrak{A} , then we know for any element $b \in \mathfrak{B}$, $\sigma_{\mathfrak{A}}(b) \subset \sigma_{\mathfrak{B}}(b)$. Then we can have more results other than it. Since the spectrum is a subset of \mathbb{C} , we need some topological properties results of \mathbb{C} .

Lemma 3. *If K is any compact subset of \mathbb{C} , then $\mathbb{C} \setminus K$ has a countable components, only one of which is unbounded. And the boundary of each component is in K .*

Proof. Let $\tilde{K} = \mathbb{C} \setminus K$, then \tilde{K} is open.

Firstly, the connected component of open set in \mathbb{C} is open.

Let U be an connected component in \tilde{K} and $x \in U$. For any point $x \in U$, Since any open neighbourhood of x is connected, and K is open, there is a open neighbourhood V of x s.t. $V \subset U$.

Secondly, \mathbb{C} has just at most countable many open sets, which are pairwise disjoint.

This result is because any open set in \mathbb{C} contains a rational point.

For any two disjoint open sets A and B in \mathbb{C} , $\partial A \cap B = \emptyset$. Thus the boundary of some component of \tilde{K} can not be contained in any component of \tilde{K} , i.e. it is contained K .

Finally, since K is bounded, there is a closed ball B containing K . But the complement of B is connected, thus there is only one component of \tilde{K} containing B . Thus the other components of \tilde{K} are bounded. \square

Remark. The bounded component of $\mathbb{C} \setminus K$ is called a hole of K .

Definition 4. If $f: A \rightarrow \mathbb{C}$, where A is a set, then the norm of f on A is defined as

$$\|f\|_A = \sup \{|f(x)| : x \in A\}$$

For a compact set $K \in \mathbb{C}$, the polynomially convex hull of K is defined as

$$\hat{K} = \{ z \in \mathbb{C} : |p(z)| \leq \|p\|_K \text{ for any polynomial } p \}$$

If $K = \hat{K}$, K is called polynomially convex.

Proposition 5. Let K be a compact subset of \mathbb{C} . Then $\mathbb{C} \setminus \hat{K}$ is the unbounded component of $\mathbb{C} \setminus K$. Therefore, K is polynomially convex if and only if $\mathbb{C} \setminus K$ is connected.

Proof. Let L be the set containing K and all bounded component of $\mathbb{C} \setminus K$. Then by the Maximal Principle, $L \subset \hat{K}$. Conversely, if $\alpha \notin L$, then $(z - \alpha)^{-1}$ is analytic in a neighbourhood of L . Therefore, there is a sequence of polynomials $\{p_n\}$ s.t. $p_n \rightarrow (z - \alpha)^{-1}$. Let $q_n = (z - \alpha)p_n$. Then $q_n \rightarrow 1$, i.e. $\|q_n - 1\| < \frac{1}{2}$ for some n . But $|q_n(\alpha) - 1| = 1$, this implies $\alpha \notin \hat{K}$, i.e. $\hat{K} \subset L$. \square

By above results, now we can provide the relationships between $\sigma_{\mathfrak{A}}(b)$ and $\sigma_{\mathfrak{B}}(b)$.

Theorem 6. If \mathfrak{A} and \mathfrak{B} are Banach algebras with same identity s.t. $\mathfrak{B} \subset \mathfrak{A}$ and $b \in \mathfrak{B}$, then

- 1) $\sigma_{\mathfrak{A}}(b) \subset \sigma_{\mathfrak{B}}(b)$ and $\partial\sigma_{\mathfrak{B}}(b) \subset \partial\sigma_{\mathfrak{A}}(b)$
- 2) $\sigma_{\mathfrak{A}}^{\circ}(b) = \sigma_{\mathfrak{B}}^{\circ}(b)$
- 3) if G is a hole of $\sigma_{\mathfrak{A}}(b)$, then $G \subset \sigma_{\mathfrak{B}}(b)$ or $G \cap \sigma_{\mathfrak{B}}(b) = \emptyset$

Proof.

For 1), let $\lambda \in \partial\sigma_{\mathfrak{B}}(b)$. Since $\sigma_{\mathfrak{A}}^{\circ}(b) \subset \sigma_{\mathfrak{B}}^{\circ}(b)$, it is sufficient to show $\lambda \in \sigma_{\mathfrak{A}}(b)$. Suppose $\lambda \notin \sigma_{\mathfrak{A}}(b)$, i.e. $(b - \lambda)$ is invertible in \mathfrak{A} . But since $\lambda \in \partial\sigma_{\mathfrak{B}}(b)$, there are $\lambda_n \in \mathbb{C} \setminus \mathfrak{B}$ with $\lambda_n \rightarrow \lambda$. Thus $(b - \lambda_n)^{-1} \in \mathfrak{B}$. But $(b - \lambda_n)^{-1} \rightarrow (b - \lambda)^{-1} \in \sigma_{\mathfrak{B}}(b)$, contradicting to $\lambda \in \sigma_{\mathfrak{A}}(b)$.

2) holds because of the result of 1) and the Maximal Principle.

For 3), let $G_1 = G \cap \sigma_{\mathfrak{B}}(b)$ and $G_2 = G \setminus \sigma_{\mathfrak{B}}(b)$. Since $\partial\sigma_{\mathfrak{B}}(b) \subset \sigma_{\mathfrak{A}}(b)$ and $G \cap \sigma_{\mathfrak{A}}(b) = \emptyset$, $G_1 = G \cap \sigma_{\mathfrak{B}}^{\circ}(b)$ is open. By the facts that G_2 is clearly open and $G = G_1 \cup G_2$ and $G_1 \cap G_2 = \emptyset$, either G_1 or G_2 is empty. \square

Then we can have some useful corollaries.

Corollary 7. Let \mathfrak{A} and \mathfrak{B} be Banach algebras with same identity s.t. $\mathfrak{B} \subset \mathfrak{A}$ and $b \in \mathfrak{B}$.

1) If $\sigma_{\mathfrak{A}}(b)$ has no holes, then $\sigma_{\mathfrak{A}}(b) = \sigma_{\mathfrak{B}}(b)$.

2) If $\sigma_{\mathfrak{B}}(b) \subset \mathbb{R}$, then $\sigma_{\mathfrak{A}}(b) = \sigma_{\mathfrak{B}}(b)$.

3) $\sigma_{\mathfrak{A}}(b) = \sigma_{\mathfrak{B}}(b)$ if and only if $\rho_{\mathfrak{A}}(b)$ is connected.

Proof. 1) is clearly true since unbounded component does not intersect $\sigma_{\mathfrak{B}}(b)$.
 2) is because $\mathbb{C} \setminus \sigma_{\mathfrak{A}}(b)$ has no holes. 3) is similar as 2). \square

3.1.3 Riesz Functional Calculus

For any polynomial p with complex coefficients,

$$p(z) = \sum_{k=0}^n \alpha_k z^k$$

we can define $p(a)$ for some $a \in \mathfrak{A}$, where \mathfrak{A} is a Banach algebra

$$p(a) = \sum_{k=0}^n \alpha_k a^k$$

Clearly, $p(a)$ is well-defined. But we can do more. If f is an analytic function on $A \subset \mathbb{C}$, then f can be approximated by a sequence of polynomials

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

Similarly, we can define $f(a)$ for $a \in \mathfrak{A}$ as

$$f(a) = \sum_{n=0}^{\infty} \alpha_n a^n$$

If the radius of convergence of this sequence is R , then it can be well-defined for $\|a\| \leq R$. By the fact that $r(a) \leq \|a\|$, for the analytic function f , if $\sigma(a) \subset A$, $f(a)$ can be well-defined.

Let $\text{Hol}(a)$ denote all functions that are analytic in a neighbourhood of $\sigma(a)$. Then there is a map from $\text{Hol}(a)$ to \mathfrak{A} defined as $f \mapsto f(a)$. Now, we can find more properties of this map. Firstly, we can give another formula of $f(a)$.

If $f \in \text{Hol}(a)$, then for any closed curve γ which encloses $\text{Hol}(a)$ and any point $z_0 \in \text{Hol}(a)$,

$$f(z_0) = \frac{1}{2\pi i} \int_{\gamma} f(z)(z - z_0)^{-1} dz$$

Therefore, replacing z_0 by $a \in \mathfrak{A}$, then we have

$$f(a) = \frac{1}{2\pi i} \int_{\gamma} f(z)(z - a)^{-1} dz$$

Clearly, by the Cauchy's Integral Formula, this definition is well-defined and is coincided with above definition. But this definition can provide us a conivient method to research the map $f \mapsto f(a)$.

Theorem 1 (Riesz Functional Calculus). *Let \mathfrak{A} be a Banach algebra and $a \in \mathfrak{A}$ and the map*

$$\begin{aligned} \rho: \quad \text{Hol}(a) &\longrightarrow \mathfrak{A} \\ f &\longmapsto f(a) = \frac{1}{2\pi i} \int_{\gamma} f(z)(z-a)^{-1} dz \end{aligned}$$

has the following properties.

- 1) ρ is an algebra homomorphism.
- 2) $\rho(1) = 1$ and $\rho(z) = a$.
- 3) If $\{f_n\} \subset \text{Hol}(a)$ and $f \in \mathfrak{A}$ with $f_n \rightarrow f$ uniformly on a compact set of $\text{Hol}(a)$, then $\rho(f_n) \rightarrow \rho(f)$ in norm.

Moreover, if any map $\tau: \text{Hol}(a) \rightarrow \mathfrak{A}$ satisfies above conditions, then $\tau = \rho$.

Proof. For 1), ρ is clearly linear. And

$$\begin{aligned} f(a)g(a) &= -\frac{1}{4\pi^2} \int_{\gamma_1} f(z)(z-a)^{-1} dz \int_{\gamma_2} g(\zeta)(\zeta-a)^{-1} d\zeta \\ &= -\frac{1}{4\pi^2} \int_{\gamma_1} \int_{\gamma_2} f(z)g(\zeta) \frac{(z-a)^{-1} - (\zeta-a)^{-1}}{\zeta-z} d\zeta dz \\ &= -\frac{1}{4\pi^2} \int_{\gamma_1} f(z) \int_{\gamma_2} \frac{g(\zeta)}{\zeta-z} d\zeta (z-a)^{-1} dz \\ &\quad + \frac{1}{4\pi^2} \int_{\gamma_2} g(\zeta) \int_{\gamma_1} \frac{f(z)}{\zeta-z} dz (\zeta-a)^{-1} d\zeta \end{aligned}$$

We can choose γ_2 to enclose γ_1 , thus

$$\int_{\gamma_1} \frac{f(z)}{\zeta-z} dz = 0, \quad \int_{\gamma_2} \frac{g(\zeta)}{\zeta-z} d\zeta = 2\pi i g(z)$$

Therefore,

$$f(a)g(a) = 2\pi i \int_{\gamma_1} f(z)g(z)(z-a)^{-1} dz = (fg)(a)$$

For 2), let $f(z) = z^k$ and $\gamma = Re^{2\pi it}$, where $R > \|a\|$ and $t \in [0, 1]$, then

$$\begin{aligned}
f(a) &= \frac{1}{2\pi i} \int_{\gamma} z^k (z - a)^{-1} dz \\
&= \frac{1}{2\pi i} \int_{\gamma} z^{k-1} \left(1 - \frac{a}{z}\right)^{-1} dz \\
&= \frac{1}{2\pi i} \int_{\gamma} z^{k-1} \sum_{n=0}^{\infty} \frac{a^n}{z^n} dz \\
&= \sum_{n=0}^{\infty} \left(\frac{1}{2\pi i} \int_{\gamma} \frac{1}{z^{n-k+1}} \right) a^n \\
&= a^k
\end{aligned}$$

For 3),

$$\begin{aligned}
&\left\| \int_{\gamma} f_n(z)(z - a)^{-1} dz - \int_{\gamma} f(z)(z - a)^{-1} dz \right\| \\
&= \left\| \int_0^1 (f_n(\gamma(t)) - f(\gamma(t)))(\gamma(t) - a)^{-1} d\gamma(t) \right\| \\
&\leq \int_0^1 |f_n(\gamma(t)) - f(\gamma(t))| \|(\gamma(t) - a)^{-1}\| d|\gamma|(t) \\
&\leq M \|\gamma\| \sup \{|f_n(z) - f(z)| : z \in \gamma(t)\}
\end{aligned}$$

where M is the bound of $\|(\gamma(t) - a)^{-1}\|$ since $t \mapsto \|(\gamma(t) - a)^{-1}\|$ is continuous on $\gamma(t)$. Therefore, by the fact that $f_n \rightarrow f$ uniformly,

$$\|f_n(a) - f(a)\| \rightarrow 0$$

Finally, the uniqueness is because any $f \in \text{Hol}(a)$ can be approximated uniformly by a sequence of polynomials. Thus, 1) and 2) means $\tau(p) = \rho(p)$ for any polynomial p , and 3) provides the fact that $\tau(f) = \rho(f)$ for any $f \in \text{Hol}(a)$. \square

Remark. we have mentioned that the integral definition is coincided with the convergent definition. In fact, by 2), this statement can be proved rigorously.

Theorem 2 (Spectral Mapping Theorem). *If $a \in \mathfrak{A}$ and $f \in \text{Hol}(a)$, then*

$$\sigma(f(a)) = f(\sigma(a))$$

Proof. Firstly, there is a $g \in \text{Hol}(a)$ s.t. for $\alpha \in \sigma(a)$, $f(z) - f(\alpha) = (z - \alpha)g(z)$, that means $f(\sigma(a)) \subset \sigma(f(a))$.

Conversely, if $\alpha \notin f(\sigma(a))$, $g(z) = (f(z) - \alpha)^{-1} \in \text{Hol}(a)$. Thus, $g(a)(f(a) - \alpha) = 1$. Therefore, $\alpha \notin \sigma(f(a))$. \square

Proposition 3. Let \mathfrak{A} be a Banach algebra and $a \in \mathfrak{A}$. $\sigma(a) = F_1 \cup F_2$, where F_1 and F_2 are disjoint nonempty closed sets. Then there is a nontrivial idempotent e , i.e. $e^2 = e$, s.t.

- 1) if $ab = ba$, then $eb = be$.
- 2) if $a_1 = ae$ and $a_2 = a(1 - e)$, then $a_1a_2 = a_2a_1 = 0$.
- 3) $\sigma(a_1) = F_1 \cup \{0\}$ and $\sigma(a_2) = F_2 \cup \{0\}$.

Proof. Since F_1 and F_2 are disjoint closed set, there are two disjoint open sets G_1 and G_2 separating F_1 and F_2 . Let f be the characteristic function of G_1 and $e = f(a)$. Thus $e^2 = e$ by $f^2 = f$.

For 1), there is a more general result, $f(a)b = bf(a)$ for any $f \in \text{Hol}(a)$. It is because by extending the fact $p(a)b = bp(a)$ for any polynomial p .

2) is clearly true.

Let $f_1(z) = zf(z)$ and $f_2(z) = z(1 - f(z))$. Then $a_j = f_j(a)$ for $j = 1, 2$. Then by the Spectral Mapping Theorem $\sigma(a_j) = f_j(\sigma(a)) = F_j \cup \{0\}$. \square

Corollary 4. Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$. If $\sigma(T) = F_1 \cup F_2$, where F_1 and F_2 are disjoint nonempty closed sets, then there are closed subspaces \mathcal{H}_1 and \mathcal{H}_2 with $\mathcal{H} = \mathcal{H}_1 \oplus \mathcal{H}_2$, s.t.

- 1) $B\mathcal{H}_j \subset \mathcal{H}_j$ for $j = 1, 2$, whenever $BT = TB$.
- 2) If $T_j = T|_{\mathcal{H}_j} : \mathcal{H}_j \rightarrow \mathcal{H}_j$, then $\sigma(T_j) = F_j$, for $j = 1, 2$.

Proof. From above proposition, e is actually the projection from \mathcal{H} to \mathcal{H}_1 . Therefore by the result in above proposition, 1) clearly holds.

For 2), since e is in fact the identity from \mathcal{H}_1 to \mathcal{H}_1 , similarly, by the Spectral Mapping Theorem, $\sigma(T_j) = F_j$. \square

3.1.4 Abelian Banach Algebras

Theorem 1 (Gelfand-Mazur Theorem). If \mathfrak{A} is a Banach algebra and a division ring, then $\mathfrak{A} = \mathbb{C}$.

Proof. It is because that for any $a \in \mathfrak{A}$, $\sigma(a) \neq \emptyset$. \square

Next, we reach the structure of an abelian Banach algebra. The structure of abelian Banach algebras can be explicit by constructing a map from an abelian Banach algebra to a continuous function space on a compact space. Firstly, we can find this compact space. Let

$$\begin{aligned}\Sigma(\mathfrak{A}) &= \{\text{all algebra homomorphism } h: \mathfrak{A} \rightarrow \mathbb{C}\} \\ \mathcal{M} &= \{\text{all maximal ideals of } \mathfrak{A}\}\end{aligned}$$

for an abelian Banach algebra \mathfrak{A} . Then we can find the relationship between $\Sigma(\mathfrak{A})$ and \mathcal{M} .

Theorem 2. Let \mathfrak{A} be an abelian Banach algebra. Define a map

$$\begin{aligned} \gamma: \quad \Sigma(\mathfrak{A}) &\longrightarrow \mathcal{M} \\ h &\longmapsto \ker h \end{aligned}$$

Then γ is a bijection.

Proof. Since $\mathfrak{A}/\ker h \cong \mathbb{C}$, $\ker h \in \mathcal{M}$, i.e. γ is well-defined.

Check: $\mathfrak{A}/M \cong \mathbb{C}$ for any $M \in \mathcal{M}$

Let $\pi: \mathfrak{A} \rightarrow \mathfrak{A}/M$. If $\pi(a)$ is not invertible, then $\pi(a\mathfrak{A})$ is a proper ideal in \mathfrak{A}/M . Thus $I = \pi^{-1}(\pi(a\mathfrak{A}))$ is a proper ideal in \mathfrak{A} and $M \subset I$. Then by the maximality of M , $I = M$, i.e. $\pi(a) = 0$. In fact, for any commutative ring, this result is true. Therefore, by Gelfand-Mazur Theorem, $\mathfrak{A}/M \cong \mathbb{C}$.

Check: γ is surjective. Let $M \in \mathcal{M}$. Define $\tilde{h}: \mathfrak{A}/M \rightarrow \mathbb{C}$ as the algebraic isomorphism. Then $h = \pi \circ \tilde{h} \in \Sigma(\mathfrak{A})$ with $\ker h = M$.

Check: γ is injective. If $\ker h = \ker h'$ for $h, h' \in \Sigma(\mathfrak{A})$, then by the **Proposition 8** in the subsection 1.4.2, $h = \alpha h'$. And since $h(1) = h'(1) = 1$, $h = h'$. \square

Then we have some properties of $h \in \Sigma(\mathfrak{A})$.

Proposition 3. Let \mathfrak{A} be an abelian Banach algebra and $h \in \Sigma(\mathfrak{A})$.

- 1) h is continuous.
- 2) $\|h\| = 1$ for $h \neq 0$.

Proof. 1) holds since $\ker h$ is maximal, thus it is closed.

Let $\lambda = h(a)$. Suppose $|\lambda| > \|a\|$. Then $1 - \frac{1}{\lambda}$ is invertible. Set $b = (1 - \frac{1}{\lambda})^{-1}$, then

$$1 = h(b(1 - \frac{1}{\lambda})) = h(b) - \frac{h(b)h(a)}{\lambda} = 0$$

Therefore, $|h(a)| \leq \|a\|$ i.e. $\|h\| \leq 1$. Since $h(1) = 1$, $\|h\| = 1$. \square

Definition 4. Let \mathfrak{A} be an abelian Banach algebra. Then $\Sigma(\mathfrak{A}) \subset \mathfrak{A}^*$ endowed with the induced wk^* -topology, is called the maximal ideal space of \mathfrak{A} .

Proposition 5. If \mathfrak{A} is an abelian Banach algebra, then $\Sigma(\mathfrak{A})$ is a compact Hausdorff space. Moreover, if $a \in \mathfrak{A}$, then

$$\sigma(a) = \Sigma(a) = \{ h(a) : h \in \Sigma(\mathfrak{A}) \}$$

Proof. Since $\Sigma(\mathfrak{A}) \subset \mathfrak{A}^*$, we just need to show $\Sigma(\mathfrak{A})$ is wk^* -closed. Let $\{h_i\}$ be a net in $\Sigma(\mathfrak{A})$ s.t. $h_i \rightarrow h$ wk^* for some h in the unit closed ball of \mathfrak{A}^* . Then for $a, b \in \mathfrak{A}$,

$$h(ab) = \lim_i h_i(ab) = \lim_i h_i(a)h_i(b) = h(a)h(b)$$

and $h(1) = \lim_i h_i(1) = 1$, thus $h \in \Sigma(\mathfrak{A})$. $\Sigma(\mathfrak{A})$ is compact.

If $h \in \Sigma(\mathfrak{A})$ and $h - h(a) \in \ker h \in \mathcal{M}$, then $h - h(a)$ is not invertible, i.e. $\Sigma(a) \subset \sigma(a)$. Conversely, if $a - \lambda$ is not invertible, $(a - \lambda)\mathfrak{A}$ is a proper ideal, which can be contained in a maximal ideal. Then $(a - \lambda) \in \ker h$ with some $h \in \Sigma(\mathfrak{A})$, $\lambda = h(a) \in \Sigma(a)$. \square

Therefore, $\Sigma(\mathfrak{A})$ is the compact space we need. Then we define the map from \mathfrak{A} to $C(\Sigma(\mathfrak{A}))$.

Theorem 6. *If \mathfrak{A} is an abelian Banach algebra, the Gelfand transform is defined as*

$$\begin{aligned} \Gamma: \mathfrak{A} &\longrightarrow C(\Sigma(\mathfrak{A})) \\ a &\longmapsto \hat{a} = \Gamma(a) \end{aligned}$$

where $\hat{a}(h) = h(a)$.

1) Γ is a continuous homomorphism.

2) $\|\Gamma\| = 1$.

3) $\ker \Gamma = \bigcap \{M : M \in \mathcal{M}\}$.

4) $\|\hat{a}\|_\infty = r(a)$.

Proof. Firstly, for 1) if $h_i \rightarrow h$ in wk^* , $\hat{a}(h_i) = h_i(a) \rightarrow h(a) = \hat{a}(h)$. Thus Γ is well-defined. And

$$\Gamma(ab)(h) = h(ab) = h(a)h(b) = \Gamma(a)(h)\Gamma(b)(h)$$

Thus Γ is a homomorphism.

For 2), since $|\hat{a}(h)| = |h(a)| \leq \|a\|$, $\|\hat{a}\|_\infty \leq \|a\|$. Then $\|\Gamma\| \leq 1$. By $\Gamma(1) = 1$, $\|\Gamma\| = 1$.

For 3), $a \in \ker \Gamma$ if and only if $h(a) = 0$ for any $h \in \Sigma(\mathfrak{A})$, i.e. $a \in \bigcap \{M : M \in \mathcal{M}\}$.

For 4), it holds since $\sigma(a) = \{h(a) : h \in \Sigma(\mathfrak{A})\}$. \square

If $a \in \mathfrak{A}$ s.t. $\overline{\{p(a) : p \text{ is any polynomial}\}} = \mathfrak{A}$, then a is called a generator of \mathfrak{A} . Clearly, this \mathfrak{A} is commutative. Then we can find an extral property of this special algebra.

Proposition 7. *If \mathfrak{A} is an abelian Banach algebra with a generator a , then there is a homeomorphism $\tau: \Sigma(\mathfrak{A}) \rightarrow \sigma(a)$ s.t. $\Gamma(p(a)) = p \circ \tau$.*

Proof. In fact, τ can be defined as $\tau(h) = h(a)$. By above mention, τ is continuous and surjective. If $\tau h_1 = \tau h_2$, then $h_1(a) = h_2(a)$. By the fact that $h_1, h_2 \in \Sigma(\mathfrak{A})$ and a is a generator of \mathfrak{A} , $h_1 = h_2$. Thus τ is a bijection. And because $\Sigma(\mathfrak{A})$ is compact, f is a closed map, i.e. f is a homeomorphism.

$$\Gamma(p(a))(h) = p(\Gamma(a))(h) = p(\Gamma(a)(h)) = p(\tau(h))$$

\square

Remark. If \mathfrak{A} is generate by a , then $\Gamma: \mathfrak{A} \rightarrow C(\sigma(a))$ can be defined as $\Gamma(p(a)) = p$.

In fact, above proposition can extends to n generators. If $\{a_i\}_{i=1}^n$ are generators of \mathfrak{A} , i.e $\overline{\{p(a_1, \dots, a_n): p \text{ is any } n \text{ variables polynomial}\}} = \mathfrak{A}$, then we have similar results as above proposition.

3.2 C^* -Algebras

Now, we have known $\mathcal{B}(\mathcal{H})$ is a Banach algebra. But there is another algebraic operation on $\mathcal{B}(\mathcal{H})$, which let $\mathcal{B}(\mathcal{H})$ be more interesting than the general Banach algebra. This operation is a map $T \rightarrow T^*$ on $\mathcal{B}(\mathcal{H})$, called an involution, and moreover, it satisfies the condition $\|T\| = \|T^*\| = \|T^*T\|^{\frac{1}{2}}$. This identity provides a strong relation between the topological structure and the algebraic structure on $\mathcal{B}(\mathcal{H})$. In fact, the topology is completely determined by the algebraic structure on $\mathcal{B}(\mathcal{H})$. In order to research this structure, we firstly define an general algebra satisfying above condition, called a C^* -algebra. By digging its topological structures and algebraic structures, we can embed it into $\mathcal{B}(\mathcal{H})$ for some Hilbert space \mathcal{H} . Therefore, any C^* -algebra can be regarded as a subalgebra of $\mathcal{B}(\mathcal{H})$.

3.2.1 Elementary Properties

Definition 1. If \mathfrak{A} is a Banach algebra, an involution is a map $a \rightarrow a^*$ from \mathfrak{A} to \mathfrak{A} satisfying for any $a, b \in \mathfrak{A}$ and any $\alpha \in \mathbb{C}$,

- 1) $(a^*)^* = a$,
- 2) $(ab)^* = b^*a^*$,
- 3) $(\alpha a + b)^* = \bar{\alpha}a^* + b^*$.

Definition 2. A C^* -algebra is a Banach algebra \mathfrak{A} with an involution s.t. for every $a \in \mathfrak{A}$,

$$\|a^*a\| = \|a\|^2$$

Then we can get some easy properties for the norm and the involution.

Proposition 3. Let \mathfrak{A} be a C^* -algebra and $a \in \mathfrak{A}$.

- 1) $\|a^*\| = \|a\|$.
- 2) $\|aa^*\| = \|a\|^2$.
- 3) $\|a\| = \sup \{\|ax\| : \|x\| \leq 1\} = \sup \{\|xa\| : \|x\| \leq 1\}$.

Proof. For 1), $\|a\|^2 = \|a^*a\| \leq \|a^*\| \|a\|$, thus $\|a\| \leq \|a^*\|$. Taking the involution, $\|a^*\| \leq \|a\|$.

$$2), \|aa^*\| = \|(a^*)^*\| = \|a^*\|^2 = \|a\|^2.$$

For 3), let α be the supremum, then $\alpha \leq \|a\|$. $a = 0$ is clearly true. For $a \neq 0$, let $x = a^*/\|a\|$. Thus $\alpha \geq \|a\|$. \square

If a C^* -algebra \mathfrak{A} is without the identity, we can use same method of the Banach algebra to extend it to a unit C^* -algebra $\tilde{\mathfrak{A}}$. The only thing we need to prove is the identity. And this can be obtained the result in above proposition. Therefore, we always assume a C^* -algebra is with the identity.

Definition 4. Let \mathfrak{A} be a C^* -algebra and $a \in \mathfrak{A}$.

- 1) a is self-adjoint if $a = a^*$.
- 2) a is normal if $aa^* = a^*a$.
- 3) a is unitary if $aa^* = a^*a = 1$.
- 4) a is a projection if a is self-adjoint and $a = a^2$.

Then we can see the algebraic structure on a C^* -algebra completely determine its norm topology.

Theorem 5. Let \mathfrak{A} be a C^* -algebra and $a \in \mathfrak{A}$. If a is self-adjoint, then

$$r(a) = \|a\|$$

Proof. Since a is self-adjoint,

$$\|a\|^2 = \|a^*a\| = \|a^2\|$$

Thus by induction, we have $\|a\|^{2n} = \|a^{2n}\|$. Then

$$r(a) = \lim_{n \rightarrow \infty} \|a^n\|^{\frac{1}{n}} = \lim_{n \rightarrow \infty} \|a^{2n}\|^{\frac{1}{2n}} = \|a\| \quad \square$$

Remark. For any $b \in \mathfrak{A}$, we know b^*b is self-adjoint,

$$r(b^*b) = \|b^*b\| = \|b\|^2$$

Thus, the norm in a C^* -algebra is completely determined by the spectral radius, which is totally an algebraic trait.

Now, we can see how the algebraic property influences the topological structure.

Proposition 6. Let $\rho: \mathfrak{A} \rightarrow \mathfrak{B}$ be a $*$ -homomorphism between two C^* -algebras.

- 1) ρ is continuous, and moreover $\|\rho(a)\| \leq \|a\|$.
- 2) If ρ is a $*$ -isomorphism, then ρ is an isometry.

Proof. 2) is the direct corollary from 1). For 1), clearly $\sigma(\rho(a)) \subset \sigma(a)$, thus

$$\|\rho(a)\|^2 = r(\rho(a^*a)) \leq r(a^*a) = \|a\|^2 \quad \square$$

Let $\text{Re } \mathfrak{A}$ denote the set of all self-adjoint elements in a C^* -algebra \mathfrak{A} . Then, for any $a \in \mathfrak{A}$, there are $x, y \in \text{Re } \mathfrak{A}$, s.t.

$$a = x + iy, \text{ where } x = \frac{a + a^*}{2}, y = \frac{a - a^*}{2i}$$

Therefore, any element in a C^* -algebra \mathfrak{A} can be combined by two self-adjoint elements. And the self-adjoint element play a important role in the algebraic structure of a C^* -algebra .

Proposition 7. *If $h: \mathfrak{A} \rightarrow \mathbb{C}$ is an algebraic homomorphis on a C^* -algebra \mathfrak{A} .*

- 1) *If $a \in \text{Re } \mathfrak{A}$, $h(a) \in \mathbb{R}$.*
- 2) *For any $a \in \mathfrak{A}$, $h(a^*) = \overline{h(a)}$.*
- 3) *$h(a^*a) \geq 0 \ \forall \ a \in \mathfrak{A}$.*
- 4) *If $u \in \mathfrak{A}$ is a unitary, then $|u| = 1$.*

Proof. For 1), let $h(a) = \alpha + i\beta$ for $\alpha, \beta \in \mathbb{R}$ and $C^*(a+it)$ be the C^* -algebra generated by $a + it$ and 1, which is abelian. Therefore, $\|h\|_{C^*(a+it)} = 1$ by

Proposition 3 in the subsection **2.1.4**. Then we have

$$\begin{aligned} |h(a + it)| &\leq \|a + it\|^2 \\ &= \|(a + it)^*(a + it)\| \\ &= \|a^2 + t^2\| \\ &\leq \|a\|^2 + t^2 \end{aligned}$$

i.e.

$$\begin{aligned} \|a\|^2 + t^2 &\geq |\alpha + i(t + \beta)|^2 \\ &= \alpha^2 + (\beta + t)^2 \\ &= \alpha^2 + \beta^2 + 2\beta t + t^2 \end{aligned}$$

Therefore, for any $t \in \mathbb{R}$, $\|a\|^2 \geq \alpha^2 + \beta^2 + 2\beta t$. Thus, $\beta = 0$.

2) and 3) and 4) is the direct results from 1). \square

Corollary 8. *If $a \in \text{Re } A$, then $\sigma(a) \subset \mathbb{R}$.*

Proof. Let $C^*(a)$ be the C^* -algebra generated by 1 and a . Thus $C^*(a)$ is abelian. Then by **Proposition 5** in the subsection **2.1.4**,

$$\sigma_{C^*(a)}(a) = \{ h(a) : h \in \Sigma(\mathfrak{A}) \} \subset \mathbb{R}$$

And by the *Corollary 7* in the subsection **2.1.2**,

$$\sigma(a) = \sigma_{C^*(a)}(a) \subset \mathbb{R}$$

□

And the spectrum of a element in a C^* -algebra has better property.

Theorem 9. *If \mathfrak{B} is a C^* -subalgebra of a C^* -algebra \mathfrak{A} and $b \in \mathfrak{B}$, then*

$$\sigma_{\mathfrak{B}}(b) = \sigma_{\mathfrak{A}}(b)$$

Proof. It suffices to show that if b is invertible in \mathfrak{A} with the inverse x , then $x \in \mathfrak{B}$. Then $(a^*a)(xx^*) = 1$. Since $a^*a \in \mathfrak{B}$ and by above corollary, we know $xx^* \in \mathfrak{B}$. But $x = (xx^*)a^*$, thus $x \in \mathfrak{B}$. □

3.2.2 Abelian C^* -Algebras

The abelian C^* -algebra is firstly an abelian Banach algebra, thus all results in the subsection **3.1.4** can be applied to it. But since the involution and the related norm provide more information, we can have better results than the general abelian Banach algebra has. Firstly, we strengthen the **Theorem 6** in subsection **2.1.4**.

Theorem 1. *If \mathfrak{A} is an abelian C^* -algebra, then the Gelfand transform $\Gamma : \mathfrak{A} \rightarrow C(\Sigma(\mathfrak{A}))$ is an isometric $*$ -isomorphism.*

Proof. Firstly, Γ is a $*$ -homomorphism, since the result in **Theorem 6** and

$$\Gamma(a^*)(h) = h(a^*) = \overline{h(a)} = \overline{\Gamma(a)}(h)$$

Now, we can easily see Γ is an isometry by 4) in the **Theorem 6**,

$$\|\hat{a}\|_{\infty}^2 = \|a^*a\|_{\infty} = r(a^*a) = \|a\|^2$$

Finally, we need to check Γ is surjective. It is because $\Gamma(\mathfrak{A})$ is a closed subalgebra of $C(\Sigma(\mathfrak{A}))$, which is closed under the complex conjugate and separates points in $\Sigma(\mathfrak{A})$. Then by the Stone-Weierstrass Theorem $\Gamma(\mathfrak{A}) = C(\Sigma(\mathfrak{A}))$ □

We know if $\mathfrak{A} = C^*(a)$ for some normal element a , then \mathfrak{A} is an abelian C^* -algebra. In fact,

$$\mathfrak{A} = \overline{\{ p(a, a^*) : p(z, \bar{z}) \text{ is a polynomial} \}}$$

Then, we can modify the result in the **Proposition 7** in subsection **2.1.4**.

Theorem 2. *Let $\mathfrak{A} = C^*(a)$ for some normal element a . Then there is a unique isometric $*$ -isomorphism $\rho: \mathfrak{A} \rightarrow C(\sigma(a))$.*

Proof. Firstly, we have a similar homeomorphism

$$\begin{aligned} \tau: \quad \Sigma(\mathfrak{A}) &\longrightarrow \sigma(a) \\ h &\longmapsto h(a) \end{aligned}$$

Then $\rho(x) = \Gamma(x) \circ \tau^{-1}$ is indeed an isometric $*$ -isomorphism by the property of Gelfand transform. And moreover, by the result of **Proposition 7**, for any $z \in \sigma(a)$, $z = h(a)$ for some $h \in \Sigma(\mathfrak{A})$

$$\begin{aligned} \rho(p(a, a^*))(z) &= \rho(p(a, a^*))(h(a)) = \Gamma(p(a, a^*)) \circ \tau^{-1}(h(a)) \\ &= \Gamma(p(a, a^*))(h) = h(p(a, a^*)) \\ &= p(h(a), h(a^*)) = p(h(a), \overline{h(a)}) \\ &= p(z, \overline{z}) \end{aligned}$$

That means that ρ maps the polynomials in \mathfrak{A} to polynomials in $\sigma(a)$. Therefore, ρ is unique. \square

By the Riesz Functional Calculus, we have the map $f \mapsto f(a)$ from $\text{Hol}(a)$ to \mathfrak{A} for a C^* -algebra \mathfrak{A} and any element $a \in \mathfrak{A}$. Now, we can extend this definition a $C(\sigma(a))$ by above ρ , but we need a is normal, then

$$\begin{aligned} \rho^{-1}: \quad C(\sigma(a)) &\longrightarrow C^*(a) \\ f &\longmapsto f(a) \end{aligned}$$

defined above is an isometric isomorphism and ρ^{-1} maps

$$\begin{aligned} 1 &\longmapsto 1 \\ z &\longmapsto a \\ \overline{z} &\longmapsto a^* \\ z^{-1} &\longmapsto a^{-1} \\ p(z, \overline{z}) &\longmapsto p(a, a^*) \end{aligned}$$

Therefore, this map is unique and it is clearly the extension of the Riesz Functional Calculus, called Continuous Functional Calculus. Like the Riesz Functional Calculus, there is also a Spectral Theorem.

Theorem 3 (Spectral Theorem). *Let \mathfrak{A} be a C^* -algebra and $a \in \mathfrak{A}$ be a normal element, then for $f \in C(\sigma(a))$,*

$$\sigma(f(a)) = f(\sigma(a))$$

Proof. For some compact space X and $f \in C(X)$, then $C(X)$ with the supremum norm is a C^* -algebra and $\sigma(f) = \text{ran } f$. Then since $f \mapsto f(a)$ is a $*$ -isomorphism,

$$\sigma(f(a)) = \sigma_{C(\sigma(a))}(f) = \text{ran } f = f(\sigma(a)) \quad \square$$

There is an important example.

Example 4. Let μ be a compactly supported, regular Borel measure on \mathbb{C} and (X, Ω, μ) be the measure space. For each $\pi \in L^\infty(\mu)$, we define the map

$$\begin{aligned} M_\phi: L^2(\mu) &\longrightarrow L^2(\mu) \\ f(z) &\longmapsto \phi(z)f(z) \end{aligned}$$

Then clearly M_ϕ is in $\mathcal{B}(L^2(\mu))$.

- 1) $(M_\phi)^* = M_{\bar{\phi}}$ and M_ϕ is a normal element in $\mathcal{B}(L^2(\mu))$.
- 2) $\phi \mapsto M_\phi$ is a $*$ -homomorphism from $L^\infty(\mu)$ to $\mathcal{B}(L^2(\mu))$.
- 3) $\|M_\phi\| = \|\phi\|_\infty$.
- 4) $\sigma(M_\phi) = \bigcap \{\overline{\phi(U)} : U \in \Omega \text{ \& } \mu(X \setminus U) = 0\}$.
- 5) If $f \in C(\sigma(M_\phi))$, then $f(M_\phi) = M_{f \circ \phi}$.

If $\phi(z) = z$, we set denote $N_\mu = M_\phi$ and in fact, $\sigma(N_\mu) = \text{supp } \mu$.

3.2.3 Positive Elements

We have known that self-adjoint elements play a important role in a C^* -algebra \mathfrak{A} . The self-adjoint element a in \mathfrak{A} is like the real number in \mathbb{C} , and the relationship between them can be revealed by the fact $a \in \text{Re } \mathfrak{A}$ if and only if $\sigma(a) \subset \mathbb{R}$. The converse is obtained by the Continuous Functional Calculus. In fact, Continuous functional calculus can provide more relation between the element in \mathfrak{A} and the element in \mathbb{C} , like positivity.

Definition 1. Let \mathfrak{A} be a C^* -algebra and $a \in \text{Re } \mathfrak{A}$. Then a is called a positive element if and only if $\sigma(a) \subset \mathbb{R}^+$, denoted by $a \geq 0$. And let \mathfrak{A}_+ be the set of all positive elements.

This definition is nature by above mention, but it may not be very explicit. So we need to show more direct equivalent definitions of positive elements.

Theorem 2. Let \mathfrak{A} be a C^* -algebra . Then the following statements are equivalent.

- 1) $a \geq 0$.
- 2) $a = b^2$ for some $b \geq 0$.
- 3) $a \in \text{Re } \mathfrak{A}$ and $\|t - a\| \leq 0$ for all $t \geq \|a\|$.
- 4) $a \in \text{Re } \mathfrak{A}$ and $\|t - a\| \leq 0$ for some $t \geq \|a\|$

Proof. All of these can be done by the functional calculus. And our goal is to find some valid functions in $C(\sigma(a))$ to complete these.

1) \Rightarrow 2) Let $f(x) = \sqrt{x}$ in $C(\sigma(a))$ and since $\sigma(a) \subset \mathbb{R}^+$, f is well-defined. Let $b = f(a)$. Then, we have $a = b^2$. And by Spectral Theorem, $\sigma(b) = \sigma(f(a)) \subset \mathbb{R}^+$.

2) \Rightarrow 3) Let $f(x) = x^2$ defined on $\sigma(b)$, then $a = f(b)$ and $\|a\| = \|f\|_\infty$. By this condition, $f(x) \geq 0$. Thus $\sigma(a) = f(\sigma(b)) \subset \mathbb{R}^+$.

3) \Rightarrow 4) It is trivial.

4) \Rightarrow 1) Let $f(x) = x$ defined on $\sigma(a) \subset \mathbb{R}$. Then this condition means

$$\|t - f\|_\infty = \|t - f(a)\| = \|t - a\| \leq t$$

for some $t \geq \|a\| = \|f\|_\infty$. Therefore, $f(x) \geq 0$ for all $x \in \sigma(a)$. Thus $\sigma(a) = f(\sigma(a)) \subset \mathbb{R}^+$. \square

Like the fact that any element in a C^* -algebra can be combined by two self-adjoint elements, any self-adjoint element can be combined by two positive elements.

Proposition 3. Let \mathfrak{A} be a C^* -algebra. If $a \in \text{Re } \mathfrak{A}$, then there are unique $u, v \in \mathbb{R}^+$, s.t.

$$a = u - v \quad \& \quad uv = vu = 0$$

Proof. Let $f(x) = \max x, 0$ and $g(x) = -\min x, 0$. Then $f, g \in C(\sigma(a))$ and $f(x) - g(x) = x$ and $f(x)g(x) = 0$. Then $u = f(a)$ and $v = g(a)$ satisfy above conditions.

If $a = u_1 - v_1$, then we can know $C^*(a, u, v, u_1, v_1)$ is an abelian C^* -algebra, thus for some compact space X , $C^*(a, u, v, u_1, v_1) \cong C(X)$. And this uniqueness can be proved in a continuous function space. \square

Corollary 4. Let \mathfrak{A} be a C^* -algebra. Then \mathfrak{A}_+ is a cone.

Proof. Let $\{a_n\} \subset \mathfrak{A}_+$ be a sequence s.t. $a_n \rightarrow a$. Then by above proposition, $\|a_n - \|a_n\|\| \leq \|a_n\|$. Taking norm limit, $\|a - \|a\|\| \leq \|a\|$, thus $a \in \mathfrak{A}_+$.

Clearly, $\alpha \mathfrak{A}_+ \subset \mathfrak{A}_+$ for any $\alpha > 0$. For $a, b \in \mathfrak{A}$, we can assume that $\|a\| \leq 1$ and $\|b\| \leq 1$, then

$$\left\| 1 - \frac{1}{2}(a + b) \right\| = \frac{1}{2} \|(1 - a) + (1 - b)\| \leq 1$$

Thus $\frac{1}{2}(a + b) \in \mathfrak{A}_+$, i.e. $a + b \in \mathfrak{A}_+$. \square

Then, we can build an order on $\text{Re } \mathfrak{A}$ by defining $a \leq b \Leftrightarrow b - a \in \mathfrak{A}_+$. And moreover, let $\mathfrak{A}_- = -\mathfrak{A}_+$, then $\mathfrak{A}_- \cap \mathfrak{A}_+ = \{0\}$. There are other properties of positivity.

Proposition 5. *Let \mathfrak{A} be a C^* -algebra .*

- 1) *If $a \geq 0$, then there is a unique $b \geq 0$ s.t. $a = b^n$.*
- 2) *If $a \in \mathfrak{A}$, then $a^*a \in \mathfrak{A}_+$.*
- 3) *If $a \leq b$ in $\text{Re } \mathfrak{A}$, then $c^*ac \leq c^*bc$ for any $c \in \mathfrak{A}$.*
- 4) *For any $a \in \text{Re } \mathfrak{A}$, $-\|a\| \leq a \leq \|a\|$ and if $a \in \mathfrak{A}$, $0 \leq a^*a \leq \|a\|^2$.*
- 5) *If $0 \leq a \leq b$, then $b^{-1} \leq a^{-1}$.*
- 6) *For any $a \in \mathfrak{A}$, we define $|a| = \sqrt{a^*a}$, then $|a| = a_+ + a_-$.*
- 7) *If $0 \leq a \leq b \in \mathfrak{A}$, then $\|a\| \leq \|b\|$.*

Proof. For 1), let $f(x) = \sqrt[n]{x}$ defined on $\sigma(a) \subset \mathbb{R}$, then $b = f(a)$ satisfying $a = b^n$.

For 2), let $b = a^*a = b_+ - b_-$ and $c = \sqrt{b_+}$ and $d = ac$. Since $d^*d = b_+^2 \in \mathfrak{A}_+$, $d^*d \in \mathfrak{A}_-$. Let $d = x + iy$, then $dd^* + d^*d = 2(x^2 + y^2) \in \mathfrak{A}_+$. Thus

$$dd^* = dd^* + d^*d - d^*d \in \mathfrak{A}_+$$

By the fact that $\sigma(dd^*) \cup \{0\} = \sigma(d^*d) \cup \{0\}$, $b_-^2 = -d^*d = 0$. Therefore, $b = a^*a \in \mathfrak{A}_+$.

For 3), let $d = \sqrt{b - a}$, then

$$c^*bc - c^*ac = c^*(b - a)c = c^*d^*dc = (dc)^*dc \in \mathfrak{A}_+$$

For 4), if $x \in \sigma(a)$, $|x| \leq \|a\|$. Therefore, $f(x) = \|a\| - x$ and $g(x) = \|a\| + x$ are positive on $\sigma(a)$.

For 5), since $0 \leq a \leq b$,

$$1 - b^{-\frac{1}{2}}ab^{-\frac{1}{2}} = b^{-\frac{1}{2}}(b - a)b^{-\frac{1}{2}} \geq 0$$

i.e. $(a^{\frac{1}{2}}b^{-\frac{1}{2}})^*(a^{\frac{1}{2}}b^{-\frac{1}{2}}) \leq 1$, therefore $\left\|a^{\frac{1}{2}}b^{-\frac{1}{2}}\right\| \leq 1$ by functional calculus as similar as 4). And thus $1 \geq (a^{\frac{1}{2}}b^{-\frac{1}{2}})(a^{\frac{1}{2}}b^{-\frac{1}{2}})^* = a^{\frac{1}{2}}b^{-1}a^{\frac{1}{2}}$. Therefore, $a^{-1} = a^{-\frac{1}{2}}1a^{-\frac{1}{2}} \geq b^{-1}$.

6) holds by the functional calculus and the uniqueness is by 1).

For 7), $\sigma(a) = \{h(a) : h \in \Sigma(\mathfrak{A})\} \subset \mathbb{R}^+$ and $\sigma(b) = \{h(b) : h \in \Sigma(\mathfrak{A})\} \subset \mathbb{R}^+$ and $\sigma(b - a) = \{h(b - a) : h \in \Sigma(\mathfrak{A})\} \subset \mathbb{R}^+$, therefore $h(b) \geq h(a) \geq 0$ for any $h \in \Sigma(\mathfrak{A})$. That means $r(b) \geq r(a)$, i.e. $\|b\| \geq \|a\|$. \square

By using 2) on above proposition, we can easily see that

Corollary 6. *Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$. If T is positive, then for any $h \in \mathcal{H}$*

$$\langle Th, h \rangle \geq 0$$

Remark. There is a $A \in \mathcal{B}(\mathcal{H})$ s.t. $T = A^*A$, thus above statement is true. In fact, the converse is also true.

3.2.4 Approximate Identities

When we research the proper ideal of an algebra, this ideal does not contain the identity. So for the ideal of a C^* -algebra, we want to find some element has similar property as the identity has in the ideal.

Definition 1. Let \mathfrak{A} be a C^* -algebra and $\{e_i\}$ be a net in \mathfrak{A} s.t.

- 1) $0 \leq e_i \leq 1$ for all i ,
- 2) $e_i \leq e_j$ for $i \leq j$,
- 3) $\lim_i ae_i = \lim_i e_i a = a$ for any $a \in \mathfrak{A}$,

Then $\{e_i\}$ is called an approximate identity for \mathfrak{A} .

Theorem 2. Every C^* -algebra \mathfrak{A} has an approximate identity.

Proof. Firstly, let $\Lambda = \{e \in \mathfrak{A}_+ : e < 1\}$. We can check Λ is indeed a direct set with respect to \leq . Define two functions as

$$\begin{aligned} f(t) &= \frac{t}{1-t}, \quad \forall t \in [0, 1), \\ g(t) &= \frac{t}{1+t} = 1 - \frac{1}{1+t}, \quad \forall t \in [0, \infty). \end{aligned}$$

In fact, $g(f(t)) = t$. Then for any $a, b \in \Lambda$, let $y = f(a) + f(b)$ and $c = g(y)$. And since $\|g\|_\infty < 1$, $c \in \Lambda$. The fact that $x = f(a) \leq y$ implies $1+x \leq 1+y$. Then $(1+x)^{-1} \geq (1+y)^{-1}$.

$$a = g(f(a)) = g(x) = 1 - (1+x)^{-1} \leq 1 - (1+y)^{-1} = c$$

Similarly, $b \leq c$. Therefore, Λ is direct.

If $a \in \mathfrak{A}_+$, let $e_n = g(na) \in \Lambda$. Define

$$h(t) = t^2(1 - g(nt)) = \frac{t^2}{1+nt} \leq \frac{t}{n}$$

Thus $h(a) = a^2(1 - g(na)) = a(1 - e_n)a$, that means

$$\|a(1 - e_n)a\| = \|h\|_\infty \leq \frac{\|a\|}{n}$$

For any $\varepsilon > 0$, we can choose a N , s.t. for $n > N$, $\|a(1 - e_n)a\| < \varepsilon$. Moreover, since for $0 \leq d \leq b \leq 1 \in \mathfrak{A}$, $c^*(1-b)c \leq c^*(1-d)c$ for any $c \in \mathfrak{A}$. Therefore,

$$\|c^*(1-b)c\| \leq \|c^*(1-d)c\|$$

And combining above mention and the fact for $0 \leq d \leq b \leq 1 \in \mathfrak{A}$,

$$\begin{aligned} \|c - dc\|^2 &\leq \|c^*(1-d)c\| \\ \|c - cd\|^2 &\leq \|c^*(1-d)c\| \end{aligned}$$

implies for $e \geq e_N$,

$$\begin{aligned}\|a - ea\|^2 &< \varepsilon \\ \|a - ae\|^2 &< \varepsilon\end{aligned}$$

Therefore,

$$\lim_i ae_i = \lim_i e_i a = a, \quad \forall a \in \mathfrak{A}$$

For arbitrary $a \in \mathfrak{A}$, a can be write as the linear combination of four positive elements. \square

Remark. The result that Λ is direct is also true for \mathfrak{A} without the identity, since $g(0) = f(0) = 0$, which means $x, y, c \in \mathfrak{A}$. And if \mathfrak{A} is separable, we can find a sequential approximate identity in its countable dense subset. Moreover, this sequential approximate identity can apply in whole \mathfrak{A} .

Then we can use the approximate identity to get some interesting results. First, we need a lemma.

Lemma 3. *If \mathfrak{A} is a C^* -algebra and $x, y \in \mathfrak{A}$, $a \in \mathfrak{A}_+$ s.t. $x^*x \leq a^\alpha$ and $y^*y \leq a^\beta$ for some positive scalars α and β with $\alpha + \beta > 1$, then the sequence $u_n = x(n^{-1} + a)^{-\frac{1}{2}}y$ converges to a $u \in \mathfrak{A}$ with $\|u\| \leq \|a^{\frac{1}{2}(\alpha+\beta-1)}\|$.*

Proof. Let $d_n m = (n^{-1} + a)^{-\frac{1}{2}} - (m^{-1} + a)^{-\frac{1}{2}}$.

$$\begin{aligned}\|u_n - u_m\|^2 &= \|xd_nmy\|^2 = \|y^*d_nmx^*xd_nmy\| \\ &\leq \left\| a^{\frac{\alpha}{2}}d_nmy \right\|^2 = \left\| a^{\frac{\alpha}{2}}d_nmy^*yd_nma^{\frac{\alpha}{2}} \right\| \\ &\leq \left\| a^{\frac{\alpha}{2}}d_nma^\beta d_nma^{\frac{\alpha}{2}} \right\| \\ &= \left\| d_nma^{\frac{\alpha+\beta}{2}} \right\|^2\end{aligned}$$

Since $f_n(t) = (n^{-1} + t)t^{\frac{\alpha+\beta}{2}}$ is an increasing positive sequence, $d_n = (n^{-1} + a)a^{\frac{\alpha+\beta}{2}}$ is an increasing positive sequence in \mathfrak{A}_+ . By Dini's Theorem, since $\sigma(a)$ is compact, there is a continuous fuction f s.t. $f_n \rightarrow f$ uniformly. Let $d = f(a)$, thus $d_n \rightarrow d$ in norm. Thus $\{u_n\}$ is Cauchy and there exists $u = \lim_{n \rightarrow \infty} u_n$.

$$\|u_n\| = \left\| x(n^{-1} + a)^{-\frac{1}{2}} \right\| \leq \left\| (n^{-1} + a)a^{\frac{\alpha+\beta}{2}} \right\| \leq \left\| a^{\frac{1}{2}(\alpha+\beta-1)} \right\| \quad \square$$

Proposition 4. *If \mathfrak{A} is a C^* -algebra and $a \in \mathfrak{A}_+$ and $x \in \mathfrak{A}$ with $x^*x \leq a$, and $0 < \alpha < \frac{1}{2}$, then there exists $u \in \mathfrak{A}$ with $\|u\| \leq \|a^{\frac{1}{2}-\alpha}\|$ and $x = ua^\alpha$.*

Proof. Let $u_n = x(n^{-1} + a)^{-\frac{1}{2}}a^{\frac{1}{2}-\alpha}$. By above lemma, $u_n \rightarrow u$ with $\|u\| \leq \|a^{\frac{1}{2}-\alpha}\|$. We use similar prove as above lemma,

$$\begin{aligned}\|x - u_n a^\alpha\|^2 &= \left\|x(1 - (n^{-1} + a)^{-\frac{1}{2}}a^{\frac{1}{2}})\right\|^2 \\ &\leq \left\|(1 - (n^{-1} + a)^{-\frac{1}{2}}a^{\frac{1}{2}})a(1 - (n^{-1} + a)^{-\frac{1}{2}}a^{\frac{1}{2}})\right\| \\ &= \left\|a^{\frac{1}{2}}(1 - (n^{-1} + a)^{-\frac{1}{2}}a^{\frac{1}{2}})\right\|^2 \\ &= \left\|a^{\frac{1}{2}} - (n^{-1} + a)^{-\frac{1}{2}}a\right\|^2\end{aligned}$$

By the Dini's Theorem, $(n^{-1} + a)^{-\frac{1}{2}}a \rightarrow a$ in norm. Therefore, $u_n a^\alpha \rightarrow a$ i.e. $a = u a^\alpha$. \square

Corollary 5. *If \mathfrak{A} is a C^* -algebra and $x \in \mathfrak{A}$ and $0 < \beta < 1$, then there is a $u \in \mathfrak{A}$ s.t.*

$$x = u |x|^\beta$$

3.2.5 Ideals and Quotients

Firstly, there are two easy results of closed ideal in a C^* -algebra.

Proposition 1. *Let \mathfrak{A} be a C^* -algebra.*

- 1) *If \mathfrak{J} is a closed left or right ideal of \mathfrak{A} , $a \in \mathfrak{J}$ with $a = a^*$, then for $f \in C(\sigma(a))$ with $f(0) = 0$, $f(a) \in \mathfrak{A}$.*
- 2) *If \mathfrak{J} is a closed ideal, then $a \in \mathfrak{J}$ implies $a^* \in \mathfrak{A}$.*

Proof. For 1), if \mathfrak{J} is proper, $0 \in \sigma(a)$. Then $f(0) = 0$ and $\sigma(a) \subset \mathbb{R}$, f can be approximated by a sequence of polynomials p_n with $p_n(0) = 0$. Therefore, $p_n(a) \in \mathfrak{J}$ and by the fact that \mathfrak{J} is closed, $f(a) \in \mathfrak{J}$.

For $a \in \mathfrak{J}$, by the corollary in above subsection, we know there is a $u \in \mathfrak{A}$ s.t. $a = u |a|^{\frac{1}{2}}$. By 1), $|a|^{\frac{1}{2}} \in \mathfrak{J}$. Therefore,

$$a^* = |a|^{\frac{1}{2}} u \in \mathfrak{J} \quad \square$$

Definition 2. If \mathfrak{A} is a C^* -algebra and \mathfrak{B} is a $*$ -subalgebra of \mathfrak{A} , then \mathfrak{B} is called hereditary if for any $b \in \mathfrak{B}_+$ and $x \in \mathfrak{A}$ with $0 \leq x \leq b$, $x \in \mathfrak{B}$.

Now, we can give more profound properties of the closed left ideals.

Theorem 3. *Let \mathfrak{A} be a C^* -algebra.*

- 1) *If \mathfrak{J} is a closed left ideal of \mathfrak{A} and $\mathfrak{B} = \mathfrak{J} \cap \mathfrak{J}^*$, then \mathfrak{B} is a hereditary subalgebra of \mathfrak{A} .*

- 2) If \mathfrak{B} is a hereditary subalgebra of \mathfrak{A} and $\mathfrak{I} = \{x \in \mathfrak{A} : x^*x \in \mathfrak{B}\}$, then \mathfrak{I} is a closed left ideal of \mathfrak{A} .
- 3) If \mathfrak{I} is a closed left ideal of \mathfrak{A} and $\mathfrak{B} = \mathfrak{I} \cap \mathfrak{I}^*$, then $\mathfrak{I} = \{x \in \mathfrak{A} : x^*x \in \mathfrak{B}\}$.
- 4) If \mathfrak{B} is a hereditary subalgebra of \mathfrak{A} and $\mathfrak{I} = \{x \in \mathfrak{A} : x^*x \in \mathfrak{B}\}$, then $\mathfrak{B} = \mathfrak{I} \cap \mathfrak{I}^*$.

Proof. For 1), clearly, \mathfrak{B} is a C^* -algebra. Let $0 \leq x \leq b$ for some $b \in \mathfrak{B}_+$. Since $x^{\frac{1}{2}}x^{\frac{1}{2}} \leq b$ there is a $u \in \mathfrak{A}$, s.t. $x^{\frac{1}{2}} = ub^{\frac{1}{3}}$. But $b^{\frac{1}{3}} \in \mathfrak{B}_+ \subset \mathfrak{I}$. Therefore, $x^{\frac{1}{2}} \in \mathfrak{I}$ and thus $x \in \mathfrak{I}$. Since x is self-adjoint, $x \in \mathfrak{B}$.

For 2), if $x \in \mathfrak{I}$ and $a \in \mathfrak{A}$, then

$$(ax)^*ax = x^*a^*ax \leq \|a\|^2 x^*x \in \mathfrak{B}.$$

Therefore, $(ax)^*ax \in \mathfrak{B}$, i.e. $ax \in \mathfrak{I}$. If $x, y \in \mathfrak{I}$, then

$$(x+y)^*x + y \leq (x+y)^*(x+y) + (x-y)^*(x-y) = 2(x^*x + y^*y) \in \mathfrak{B}$$

Thus, \mathfrak{I} is a closed left ideal.

For 3), if $x \in \mathfrak{A}$ and $x^*x \in \mathfrak{B}$, $x^*x \in \mathfrak{I}$. Thus $|x|^{\frac{1}{2}} \in \mathfrak{I}$ since \mathfrak{I} is a closed left ideal. Therefore, $x = u|x|^{\frac{1}{2}} \in \mathfrak{I}$. The converse is clearly true.

For 4), if $x \in \mathfrak{I}_+$, then $x^2 \in \mathfrak{B}$. So $x = \sqrt{x^2} \in \mathfrak{B}_+$. Conversely, $x \in \mathfrak{B}_+$ means $(x^{\frac{1}{2}})^*(x^{\frac{1}{2}}) = x \in \mathfrak{B}$, $x^{\frac{1}{2}} \in \mathfrak{I}_+$. Therefore, $\mathfrak{B}_+ = \mathfrak{I}_+$. Thus, $\mathfrak{B} = \mathfrak{I} \cap \mathfrak{I}^*$. \square

Theorem 4. If \mathfrak{I} is a closed ideal of a C^* -algebra \mathfrak{A} , then the quotient algebra $\mathfrak{A}/\mathfrak{I}$ with the induced norm and the induced involution, i.e. $a + \mathfrak{I}^* = a^* + \mathfrak{I}$ is also a C^* -algebra and the norm can be

$$\|a + \mathfrak{I}\| = \inf \{ \|a - ax\| : x \in \mathfrak{I}_+, \|x\| \leq 1 \}$$

Proof. Firstly, let $\{e_i\}$ be an approximate identity for \mathfrak{I} . Since $0 \leq e_i \leq 1$, if $a \in \mathfrak{A}$ and $y \in \mathfrak{I}$, $\|(a+y)(1-e_i)\| \leq \|a+y\|$.

$$\begin{aligned} \|a+y\| &\geq \liminf_i \|(a+y)(1-e_i)\| \\ &= \liminf_i \|(a-ae_i) + (y-ye_i)\| \\ &= \liminf_i \|(a-ae_i)\| \end{aligned}$$

The converse is clearly true.

The involution defined on the quotient algebra is well-defined and satisfies the conditions since \mathfrak{I} is self-adjoint. We just need to show the induced norm satisfies the C^* -identity. For $a \in \mathfrak{A}$, we have $\|a^* + \mathfrak{I}\| = \|a + \mathfrak{I}\|$, and

$$\begin{aligned}\|a + \mathfrak{I}\|^2 &= \|(a^* + \mathfrak{I})(a + \mathfrak{I})\| \\ &\leq \|a^* + \mathfrak{I}\| \|a + \mathfrak{I}\| \\ &= \|a + \mathfrak{I}\|^2\end{aligned}$$

Conversely,

$$\begin{aligned}\|a + \mathfrak{I}\|^2 &= \inf \{ \|a - ax\|^2 : x \in \mathfrak{I}_+, \|x\| \leq 1 \} \\ &= \inf \{ \|(1-x)a^*a(1-x)\| : x \in \mathfrak{I}_+, \|x\| \leq 1 \} \\ &\leq \inf \{ \|a^*a(1-x)\| : x \in \mathfrak{I}_+, \|x\| \leq 1 \} \\ &= \|a^*a + \mathfrak{I}\|\end{aligned}$$

Therefore, $\mathfrak{A}/\mathfrak{I}$ is indeed a C^* -algebra. \square

Then it can provide a C^* -algebra with some similar results as general rings have. These results can also show us how the algebraic structure in a C^* -algebra affect the topological structure. In my opinion, for researching a C^* -algebra, researching the algebraic structure may be more important.

Corollary 5. *Let \mathfrak{A} and \mathfrak{C} be two C^* -algebras.*

- 1) *If $\rho: \mathfrak{A} \rightarrow \mathfrak{C}$ is a $*$ -homomorphism, then $\text{ran } \rho$ is a C^* -algebra and the induced map $\tilde{\rho}: \mathfrak{A}/\ker \rho \rightarrow \text{ran } \rho$ is an $*$ -isomorphism.*
- 2) *If \mathfrak{I} is a closed ideal of \mathfrak{A} and \mathfrak{B} is a subalgebra of \mathfrak{A} , then there is a $*$ -isomorphism*

$$\mathfrak{B}/(\mathfrak{B} \cap \mathfrak{I}) \cong (\mathfrak{B} + \mathfrak{I})/\mathfrak{I}$$

Proof. For 1), we just need to show $\text{ran } \rho$ is closed.

Since $\tilde{\rho}$ is a $*$ -monomorphism from $\mathfrak{A}/\ker \rho$ to \mathfrak{C} , by the result in **Proposition 6** in the subsection **2.1.1**, ρ is an isometry. Thus $\text{ran } \rho = \rho(\mathfrak{A})$ is closed.

For 2), there is a commutative graph like

$$\begin{array}{ccc}\mathfrak{B} & \xrightarrow{i} & \mathfrak{A} \\ \downarrow Q & \swarrow \pi & \\ \mathfrak{A}/\mathfrak{I} & & \end{array}$$

$Q = \pi \circ i: \mathfrak{B} \rightarrow \mathfrak{B}/\mathfrak{I}$ and thus $\pi^{-1}(Q(\mathfrak{B})) = \mathfrak{B} + \mathfrak{I}$. By restricting π on $\mathfrak{B} + \mathfrak{I}$, then we have

$$(\mathfrak{B} + \mathfrak{I})/\mathfrak{I} \cong Q(\mathfrak{B}) \cong \mathfrak{B}/(\mathfrak{B} \cap \mathfrak{I}) \quad \square$$

3.2.6 Positive Functionals and GNS Construction

We have known that $\mathcal{B}(\mathcal{H})$ is a C^* -algebra. In fact, defining the general C^* -algebra is to research the operator algebra. Moreover, by the properties, we can see any C^* -algebra is a $*$ -subalgebra of $\mathcal{B}(\mathcal{H})$ for some \mathcal{H} . To prove that, we just need to find a faithful representation of the fixed C^* -algebra. Therefore, the main idea is to construct a representation (π, \mathcal{H}) . How to construct it is a question. But, fortunately, some special functionals can provide us a method.

Definition 1. Let \mathfrak{A} be C^* -algebra and ϕ is a linear functional on \mathfrak{A} . ϕ is called positive if for any $a \in \mathfrak{A}_+$, $\phi(a) \leq 0$. A positive linear functional ϕ is called a state if $\phi(1) = 1$.

Remark. Let $\mathcal{S}_{\mathfrak{A}}$ denote the set of all states on \mathfrak{A} .

Then there are some properties of positive functionals.

Proposition 2. Let \mathfrak{A} be a C^* -algebra and ϕ be a positive functional.

1) For any $x, y \in \mathfrak{A}$, then

$$|\phi(y^*x)| \leq \phi(x^*x)\phi(y^*y)$$

2) ϕ is bounded and if $\{e_i\}$ is an approximate identity of \mathfrak{A} , then

$$\|\phi\| = \lim_i \phi(e_i)$$

Proof. For 1), it can easily check that $\langle x, y \rangle = \phi(y^*x)$ is a semi-inner product, thus by the CBS inequality, above inequality is true.

Assume \mathfrak{A} without identity. If ϕ is unbounded, then there is a sequence $\{a_k\} \in \mathfrak{A}_+$ with $\|a_k\| \leq 1$ s.t. $\phi(a_k) > 2^k$. Let $a = \sum_{k=1}^{\infty} 2^{-k} a_k$. Then

$$\phi(a) \geq \phi\left(\sum_{k=1}^n 2^{-k} a_k\right) > n$$

which is a contradiction.

$$\alpha = \sup \{ \phi(a) : a \in \mathfrak{A}_+, \|a\| \leq 1 \} < \infty$$

Since any element in \mathfrak{A} can be a linear combination of four positive elements, $\|\phi\| \leq 4\alpha$. Therefore, ϕ is bounded.

Let $\beta = \lim_i \phi(e_i)$. Clearly, $\beta \leq \|\phi\|$. And since for $a \in \mathfrak{A}$ with $\|a\| \leq 1$ then $0 \leq a^*a \leq 1$

$$|\phi(a)|^2 = \lim_i |\phi(e_i a)|^2 \leq \lim_i \phi(e_i) \phi(a^* a) \leq \beta \|\phi\|$$

$$\|\phi\|^2 \leq \beta \|\phi\|$$

i.e. $\|\phi\| \leq \beta$. □

Remark. For 2), if \mathfrak{A} has an identity, then $\|\phi\| = \phi(1)$. It is because for $a \in \mathfrak{A}$ with $\|a\| \leq 1$

$$|\phi(a)|^2 \leq \phi(a^*a)\phi(1) \leq \phi(1)^2$$

i.e. $|\phi(a)| \leq \phi(1)$.

In fact, the converse of 2) in above proposition is also true.

Proposition 3. *Let \mathfrak{A} be C^* -algebra and ϕ is a bounded linear functional with $\|\phi\| = \phi(1)$, then $\phi \leq 0$.*

Proof. If $\mathfrak{A} = C(X)$ for some compact space X , ϕ is a measure μ with $\mu(X) = \|\mu\|$, then $\mu \geq 0$, i.e. ϕ is positive. Then for any \mathfrak{A} , if $a \in \mathfrak{A}_+$, $\mathfrak{B} = C^*(a) \cong C(\sigma(a))$. Then $\phi|_{\mathfrak{B}}(1) \leq \|\phi\| = \phi(1) = \phi|_{\mathfrak{B}}(1)$, thus $\phi|_{\mathfrak{B}}(a) = \phi(a) \geq 0$. \square

Using above proposition and the Hahn-Banach Theorem, we can get the corollary.

Corollary 4. *If \mathfrak{A} is a C^* -algebra and \mathfrak{B} is a C^* -subalgebra of \mathfrak{A} , then every state on \mathfrak{B} can extend to \mathfrak{A} .*

Now, we can construct the a representation of a C^* -algebra. A representation (π, \mathcal{H}) of a C^* -algebra \mathfrak{A} is called cyclic if there is a unit cyclic vector e s.t. $\pi(\mathfrak{A})e = \mathcal{H}$.

Theorem 5 (Gelfand-Naimark-Segal Construction). *Let \mathfrak{A} be a C^* -algebra and $\mathcal{S}_{\mathfrak{A}}$ be the coincided state space.*

- 1) *If $\phi \in \mathcal{S}_{\mathfrak{A}}$, then there is a cyclic representation $(\pi_{\phi}, \mathcal{H}_{\phi})$ with the unit cyclic vector e_{ϕ} s.t.*

$$\phi(a) = \langle \pi_{\phi}(a)e_{\phi}, e_{\phi} \rangle, \quad \forall a \in \mathfrak{A}$$

- 2) *If (π, \mathcal{H}) is a cyclic representation with the unit cyclic vector e , then there is a $\phi \in \mathcal{S}_{\mathfrak{A}}$ defined as*

$$\phi(a) = \langle \pi(a)e, e \rangle, \quad \text{for } a \in \mathfrak{A}$$

And for the $(\pi_{\phi}, \mathcal{H}_{\phi})$ defined as above mention, $\pi_{\phi} \cong \pi$.

Proof. For 1), the prove can be completed by several nature steps.

- 1) Constructing semi-inner product: Define a semi-inner product $\langle \cdot, \cdot \rangle$ on \mathfrak{A} like

$$\langle x, y \rangle = \phi(y^*x), \quad \text{for } x, y \in \mathfrak{A}$$

- 2) Constructing \mathcal{H}_ϕ : $\langle \cdot, \cdot \rangle$ is just a semi-inner product on \mathfrak{A} , thus we need to make it be nondegenerate. The nature method is constructing a equilateral relation or a closed ideal, and then inducing a quotient space. Naturally, we define

$$\mathfrak{B} = \{ x \in \mathfrak{A} : \phi(x) = 0 \}$$

Clearly, \mathfrak{B} is a hereditary subalgebra. Then by the **Theorem 3** in the subsection **2.2.5**, \mathfrak{B} can induced a closed left ideal

$$\mathfrak{J} = \{ x \in \mathfrak{A} : x^*x \in \mathfrak{B} \} = \{ x \in \mathfrak{A} : \phi(x^*x) = 0 \}$$

And define the inner product on $\mathfrak{A}/\mathfrak{J}$ as

$$\langle x + \mathfrak{J}, y + \mathfrak{J} \rangle = \phi(y^*x)$$

Therefore, we can easily check that $\mathfrak{A}/\mathfrak{J}$ is a inner product space. Then let \mathcal{H}_ϕ be the completion of $\mathfrak{A}/\mathfrak{J}$.

- 3) Constructing π_ϕ : Firstly, let $\pi_\phi(a)$ be defined on $\mathfrak{A}/\mathfrak{J}$ for $a \in \mathfrak{A}$.

$$\begin{aligned} \pi_\phi(a): \quad \mathfrak{A}/\mathfrak{J} &\longrightarrow \mathfrak{A}/\mathfrak{J} \\ x + \mathfrak{J} &\longmapsto ax + \mathfrak{J} \end{aligned}$$

But since

$$\begin{aligned} \|ax + \mathfrak{J}\|^2 &= \langle ax + \mathfrak{J}, ax + \mathfrak{J} \rangle = \phi(ax^*ax) \\ &\leq \|a\|^2 \phi(x^*x) = \|a\|^2 \|x + \mathfrak{J}\|^2 \end{aligned}$$

$\|\pi_\phi(a)\| \leq \|a\|$. Therefore, $\pi_\phi(a)$ can extend to \mathcal{H}_ϕ for any $a \in \mathfrak{A}$. $\pi_\phi: \mathfrak{A} \rightarrow \mathcal{B}(\mathcal{H}_\phi)$ is a representation.

- 4) Check the conditions: Let $e_\phi = 1 + \mathfrak{J}$. Then

$$\pi_\phi(\mathfrak{A})e_\phi = \{ a + \mathfrak{J} : a \in \mathfrak{A} \} = \mathfrak{A}/\mathfrak{J}$$

Therefore, π_ϕ with e_ϕ is indeed a cyclic representation of \mathfrak{A} . And clearly,

$$\langle \pi_\phi(a)e_\phi, e_\phi \rangle = \phi(a)$$

For 2), we just need to construct a unitary from \mathcal{H}_ϕ to \mathcal{H} . By observation,

$$\langle \pi_\phi(a)e_\phi, e_\phi \rangle = \langle \pi(a)e, e \rangle$$

and the facts that $\mathcal{H}_\phi = \overline{\pi_\phi(\mathfrak{A})e_\phi}$ and $\mathcal{H} = \overline{\pi(\mathfrak{A})e}$, we can find the unitary U . Firstly, let U be defined on $\pi_\phi(\mathfrak{A})e_\phi$

$$\begin{aligned} U: \quad \pi_\phi(\mathfrak{A})e_\phi &\longrightarrow \pi(\mathfrak{A})e \\ \pi_\phi(a)e_\phi &\longmapsto \pi(a)e \end{aligned}$$

Since

$$\|\pi(a)e\|^2 = \langle \pi(a)e, \pi(a)e \rangle = \phi(a^*a) = \|\pi_\phi(a)e_\phi\|^2$$

U can extend to an unitary from \mathcal{H}_ϕ to \mathcal{H} . And

$$U\pi_\phi(a)(\pi_\phi(x)e_\phi) = U\pi_\phi(ax)e_\phi = \pi(ax)e = \pi(a)(\pi(x)e) = \pi(a)U(\pi_\phi(x)e_\phi)$$

Then $U\pi_\phi U^* = \pi$, $\pi_\phi \cong \pi$. \square

Remark. In fact, in above theorem, if ϕ is a positive functional, the results is also true.

By above theorem, our faithful representation of a C^* -algebra can not be constructed by a state, since π_ϕ is not injective. But we may choose enough many states to construct a faithful representation. Therefore, we need more properties of states.

Proposition 6. *If \mathfrak{A} is C^* -algebra and $a \in \mathfrak{A}$ is self-adjoint, let $\alpha = \min \sigma(a)$ and $\beta = \max \sigma(a)$, then*

$$[\alpha, \beta] = \{ \phi(a) : \phi \in \mathcal{S}_{\mathfrak{A}} \}$$

Proof. Let $\mathfrak{B} = C^*(a)$. Then $\mathfrak{B} = \{f(a) : f \in C(\sigma(a))\}$. For $\phi \in \mathcal{S}_{\mathfrak{A}}$ and $\phi_0 = \phi|_{\mathfrak{B}}$, then there is a measure μ s.t.

$$\phi(f(a)) = \phi_0(f(a)) = \int_{\sigma(a)} f d\mu, \quad \forall f \in C(\sigma)$$

In particular, $\phi(a) = \int_{\sigma(a)} t d\mu \in [\alpha, \beta]$.

Conversely, if $\alpha \leq t_0 \leq \beta$, define $\phi_0 \in \mathcal{S}_{\mathfrak{B}}$ as

$$\phi_0(f(a)) = \frac{t_0 - \alpha}{\beta - \alpha} f(\alpha) + \frac{\beta - t_0}{\beta - \alpha} f(\beta)$$

Then ϕ_0 can extend to \mathfrak{A} , and $t_0 = \phi(a)$. \square

By this proposition, we can get an important corollary.

Corollary 7. *If \mathfrak{A} is a C^* -algebra and $a \in \mathfrak{A}$ and \mathcal{S} is wk^* -dense subset of $\mathcal{S}_{\mathfrak{A}}$, then*

$$\|a\|^2 = \sup \{ \phi(a^*a) : \phi \in \mathcal{S} \}$$

Remark. In fact, we can easily check that $\mathcal{S}_{\mathfrak{A}} \subset \mathfrak{A}^*$ is a wk^* -compact convex subset. And this corollary is equivalent to saying that \mathcal{S} can separate the points in \mathfrak{A}_+ .

Now, using this corollary, we can finally construct a faithful representation.

Theorem 8 (Gelfand-Naimark Theorem). *Every C^* -algebra \mathfrak{A} has a faithful representation (π, \mathcal{H}) . Moreover, \mathcal{H} is separable if and only if there are countable number of states on \mathfrak{A} that can separates points in \mathfrak{A}_+ , each of which defines a separable representation. In particular, each separable C^* -algebra has a faithful, separable representation.*

Proof. Let \mathcal{S} be wk^* -dense subset of $\mathcal{S}_{\mathfrak{A}}$. Define

$$\mathcal{H} = \oplus \{ \mathcal{H}_\phi : \phi \in \mathcal{S} \}$$

$$\pi = \oplus \{ \pi_\phi : \phi \in \mathcal{S} \}$$

Then we can see

$$\begin{aligned} \|a\|^2 &= \sup \{ \phi(a^*a) : \phi \in \mathcal{S} \} \\ &= \sup \{ \langle \pi_\phi(a)e_\phi, \pi_\phi(a)e_\phi \rangle : \phi \in \mathcal{S} \} \\ &= \sup \{ \|\pi_\phi(a)e_\phi\|^2 : \phi \in \mathcal{S} \} \\ &= \|\pi(a)e\|^2 \\ &\leq \|\pi(a)\|^2 \end{aligned}$$

And in the GNS Construction, we have seen $\|\pi_\phi(a)\| \leq \|a\|$ for any $\phi \in \mathcal{S}_{\mathfrak{A}}$. Therefore,

$$\|a\| \geq \|\pi(a)\|$$

Thus $\|\pi(a)\| = \|a\|$ i.e. (π, \mathcal{H}) is indeed a faithful representation.

If \mathcal{H} is separable, let $\{e_n\}$ be the dense subset of $\{h \in \mathcal{H} : \|h\| = 1\}$. Define $\mathcal{H}_n = \overline{\pi(\mathfrak{A})e_n}$ and $\pi_n = \pi|_{\mathcal{H}_n}$. Therefore, in above construction,

$$\{(\pi_\phi, \mathcal{H}_\phi, e_\phi) : \phi \in \mathcal{S}\}$$

can be replaced by

$$\{(\pi_n, \mathcal{H}_n, e_n) : n \in \mathbb{N}\}$$

Then for $a \in \mathfrak{A}_+$, there exists $b \in \mathfrak{A}$ s.t. $a = b^*b$, if for $n \in \mathbb{N}$

$$0 = \phi_n(a) = \langle \pi_n(b^*b)e_n, e_n \rangle = \|\pi(b)e_n\|^2 \quad (*)$$

But $\|b\| = \sup_n \|\pi(b)e_n\|$, that implies $a = 0$.

Conversely, if $\{\phi_n\}$ is these states, we can use the GNS Construction to get $\{(\pi_n, \mathcal{H}_n, e_n) : n \in \mathbb{N}\}$. Similarly, we can use direct sum to get the representation (π, \mathcal{H}) . Then by using the (*), we can know this representation is definitely faithful.

Finally, if \mathfrak{A} is separable, then the closed unit ball in \mathfrak{A}^* is wk^* -compactly metrizable. Therefore, there is a countable wk^* -dense subset of $\mathcal{S}_{\mathfrak{A}}$. Thus by above corollary, this subset can separate the points in \mathfrak{A}_+ \square

Remark. Therefore, any C^* -algebra can be isometrically imbedded in a operator algebra on some Hilbert space. Researching a abstract C^* -algebra is actually researching the C^* -subalgebra of an operator algebra on some Hilbert space.

Now, we can see the importance of states or positive funtionals. Thus, we want to find the relation between general linear functionals and the positive linear functionals. Firstly, we observe that for a positive functional ϕ , $\phi(a^*) = \overline{\phi(a)}$ for all a .

Definition 9. If \mathfrak{A} is a C^* -algebra and L is a linear functional, L is called self-adjoint if

$$L(a^*) = \overline{L(a)}, \quad \forall a \in \mathfrak{A}$$

Or equivalently, $L(a) \in \mathbb{R}$ for all $a \in \text{Re } \mathfrak{A}$.

The bounded self-adjoint linear functional in \mathfrak{A}^* is like the self-adjoint element in \mathfrak{A} . So by the fact that any self-adjoint element can be a linear combination of two positive element, we want to find similar result of bounded linear functionals. There is a clearly lemma of bounded self-adjoint linear functional.

Lemma 10. *If \mathfrak{A} is a C^* -algebra and L is a bounded self-adjoint linear functional, then*

$$\|L\| = \sup \{ L(a) : a \in \text{Re } \mathfrak{A}, \|a\| \leq 1 \}$$

Theorem 11 (Jordan Decomposition). *If \mathfrak{A} is a C^* -algebra and L is a bounded self-adjoint linear functional, then there are positive linear functionals ϕ_+ and ϕ_- s.t.*

$$L = \phi_+ - \phi_- \quad \text{and} \quad \|L\| = \|\phi_+\| + \|\phi_-\|$$

Proof. Let Ω be the set of all bounded self-adjoint linear functional in the unit ball of \mathfrak{A}^* and $\tilde{\Omega}$ be the wk^* -closed convex hull of $\mathcal{S}_{\mathfrak{A}} \cup (-\mathcal{S}_{\mathfrak{A}})$.

1) Claim: $\Omega = \tilde{\Omega}$

Clearly, $\tilde{\Omega} \subset \Omega$. If there is a $L_0 \in \Omega \setminus \tilde{\Omega}$, then there exists a $x_0 \in \mathfrak{A}$ separating L_0 and $\tilde{\Omega}$, and let $a_0 = \text{Re } x_0$, we have

$$L(a_0) \leq \alpha < L_0(a_0), \quad \forall L \in \tilde{\Omega}$$

That means

$$\|a_0\| = \sup \{ |\phi(a_0)| : \phi \in \mathcal{S}_{\mathfrak{A}} \} \leq \alpha < L_0(a_0)$$

but $\|L_0\| \leq 1$, which is a contradiction.

2) Claim: $\Omega = \{s\phi - t\psi : \phi, \psi \in \mathcal{S}_{\mathfrak{A}}, s + t = 1 \text{ \& } s, t \geq 0\}$

By the fact that $(s, t, \phi, \psi) \rightarrow s\phi - t\psi$ is continuous and $\mathcal{S}_{\mathfrak{A}}$ is wk^* -compact, above claim is clearly true.

3) Existence: Therefore, for any self-adjoint linear functional L with $\|L\| = 1$, there are two $\phi, \psi \in \mathcal{S}_{\mathfrak{A}}$, and positive real numbers s, t with $s + t = 1$ s.t. $L = s\phi - t\psi$. Then let $\phi_+ = s\phi$ and $\phi_- = t\psi$, then $L = \phi_+ - \phi_-$ and

$$\|\phi_+\| + \|\phi_-\| = s + t = 1 = \|L\|$$

For general L , just let $L/\|L\|$ for $L \neq 0$, the statements are true. \square

In fact, for any $f \in \mathfrak{A}^*$, $f = \phi + i\psi$ for some self-adjoint linear functional ϕ, ψ , where

$$\phi(a) = \frac{f(a) + \overline{f(a^*)}}{2} \text{ \& } \psi(a) = \frac{f(a) - \overline{f(a^*)}}{2i}$$

By combining this and the Jordan Decomposition Theorem, we can have the following corollaries.

Corollary 12. *Let \mathfrak{A} be a C^* -algebra and $\pi: \mathfrak{A} \rightarrow \mathcal{B}(\mathcal{H})$ be the faithful representation constructed in the Gelfand-Naimark Theorem.*

- 1) *Any bounded linear functional can be a linear combination of four positive linear functionals.*
- 2) *If $L \in \mathfrak{A}^*$, then there are $g, h \in \mathcal{H}$, s.t. for any $a \in \mathfrak{A}$*

$$L(a) = \langle \pi(a)g, h \rangle$$

3.2.7 Representations

Definition 1. Let $\pi: \mathfrak{A} \rightarrow \mathcal{B}(\mathcal{H})$ be a representation of a C^* -algebra \mathfrak{A} .

- 1) If there is no invariant subspace for $\pi(\mathfrak{A})$, π is called algebraically irreducible. If there is no invariant closed subspace, π is called topologically irreducible.
- 2) If $\overline{\pi(\mathfrak{A})e} = \mathcal{H}$ for some $e \in \mathcal{H}$, π is called cyclic and e is called a cyclic vector for π .
- 3) If $\overline{\pi(\mathfrak{A})\mathcal{H}} = \mathcal{H}$, π is called non-degenerate.

Remark. Since π is $*$ -homomorphism and \mathfrak{A} is $*$ -closed, the definition of irreducibility of π is valid. In fact, if π is algebraically irreducible, then π is clearly topologically irreducible. But the converse is also true. Then we have,

π irreducible $\Rightarrow \pi$ cyclic $\Rightarrow \pi$ is non-degenerate

Theorem 2. *Every representation of a C^* -algebra is equivalent to the direct sum of cyclic representations.*

Proof. Let $\pi: \mathfrak{A} \rightarrow \mathcal{B}(\mathcal{H})$ be a representation and $e \in \mathcal{H}$ be any unit element. Then $\mathcal{M} = \overline{\pi(\mathfrak{A})e}$ is clearly a invariant space for π . If $\mathcal{M} \neq \mathcal{H}$, then $\mathcal{H} = \mathcal{M} \oplus \mathcal{M}^\perp$. We have another unit vector $e' \in \mathcal{M}^\perp$, then let $\mathcal{M}' = \overline{\pi(\mathfrak{A})e'}$ and $\mathcal{M}' \perp \mathcal{M}$. If $\mathcal{M}' \neq \mathcal{M}^\perp$, we can continue above process again. By induction and the Zorn's Lemma, the theorem can be obtained. \square

Theorem 3. *If \mathfrak{A} is a C^* -algebra and \mathfrak{I} is an ideal of \mathfrak{A} , then every representation $\rho: \mathfrak{I} \rightarrow \mathcal{B}(\mathcal{H})$ can be extended to a representation $\tilde{\rho}: \mathfrak{A} \rightarrow \mathcal{B}(\mathcal{H})$. Moreover, if ρ is non-degenerate, then $\tilde{\rho}$ is unique.*

Proof. Firstly, assume ρ is non-degenerate and by above theorem we can also assume ρ is cyclic with the cyclic vector e . Then for any $a \in \mathfrak{A}$, we define

$$\begin{aligned} \rho(\tilde{a}): \quad \rho(\mathfrak{I})e &\longrightarrow \mathcal{H} \\ \rho(x)e &\longmapsto \rho(ax)e \end{aligned}$$

And since for an approximate identity $\{e_i\}$ for \mathfrak{I} , we have

$$\begin{aligned} \|\rho(ax)e\| &= \lim_i \|\rho(ae_i x)e\| \\ &\leq \lim_i \|\rho(ae_i)\rho(x)e\| \\ &\leq \sup_i \|ae_i\| \|\rho(x)e\| \\ &\leq \|a\| \|\rho(x)e\| \end{aligned}$$

i.e. $\|\tilde{\rho}(a)\| \leq \|a\|$, $\tilde{\rho}(a)$ can extend to \mathcal{H} . Then, we can check that $\tilde{\rho}: \mathfrak{A} \rightarrow \mathcal{B}(\mathcal{H})$ is indeed a representation.

If $T \in \mathcal{B}(\mathcal{H})$ s.t. for any $x \in \mathfrak{I}$, $T\rho(x) = \rho(ax)$, then

$$T\rho(x)h = \tilde{\rho}\rho(x)h, \quad \forall x \in \mathfrak{I} \text{ \& \& } \forall h \in \mathcal{H}$$

Since ρ is non-degenerate, $T = \tilde{\rho}(a)$.

Now, if ρ is degenerate, let $\mathcal{H}_0 = \overline{\rho(\mathfrak{I})}$, and $\rho_0(x) = \rho(x)|_{\mathcal{H}_0}$, then $\rho_0: \mathfrak{I} \rightarrow \mathcal{B}(\mathcal{H}_0)$ is a non-degenerate representation. Therefore, it can extend to $\tilde{\rho}_0: \mathfrak{A} \rightarrow \mathcal{B}(\mathcal{H}_0)$. Let $\kappa: \mathfrak{A} \rightarrow \mathcal{B}(\mathcal{H}_0^\perp)$ be any representation. Then $\tilde{\rho} = \tilde{\rho}_0 \oplus \kappa$ is the extension. \square

Corollary 4. *If \mathfrak{A} is a C^* -algebra and \mathfrak{I} is an ideal of \mathfrak{A} , and (ρ, \mathcal{H}) and (κ, \mathcal{K}) are two representations of \mathfrak{A} s.t.*

- 1) $\rho|_{\mathfrak{I}}$ is a non-degenerate representation of \mathfrak{I} .
- 2) $\rho|_{\mathfrak{I}} \cong \kappa|_{\mathfrak{I}}$.

then $\rho \cong \kappa$

We know the irreducible representation play a important role in the representation theory. Therefore, we want to find the necessary and sufficient conditions make a representation be irreducible.

Definition 5. If $\mathfrak{A} \subset \mathcal{B}(\mathcal{H})$ is a C^* -subalgebra, the commutant of \mathfrak{A}

$$\mathfrak{A}' = \{ T \in \mathcal{B}(\mathcal{H}) : TS = ST, \forall S \in \mathfrak{A} \}$$

Remark. Clearly, \mathfrak{A}' is a C^* -algebra.

Theorem 6. Let \mathfrak{A} be a C^* -algebra and (ρ, \mathcal{H}) be a representation, the following statements are equivalent. Then ρ is irreducible if and only if $\rho(\mathfrak{A})' = \mathbb{C}$

Proof. If ρ is irreducible and $\mathbb{C} \neq \rho(\mathfrak{A})'$, then there is a $T \in \rho(\mathfrak{A})' \setminus \mathbb{C}$. Therefore, there is a proper projection $P \in \rho(\mathfrak{A})' \setminus \mathbb{C}$. Then $P\mathcal{H}$ is a proper invariant space for ρ , which is a contradiction. The converse is clearly true by using similar prove as above. \square

Remark. The existence of P is because that $\rho(\mathfrak{A})'$ is a von Neumann algebra.

Firstly, by the GNS construction, the representations generated by positive linear functionals are important. So we need more properties of these representations.

Proposition 7. Let \mathfrak{A} be a C^* -algebra and ψ and ϕ be two positive linear functionals on \mathfrak{A} and (π, \mathcal{H}, e) be the representation generated by ϕ . Then $\psi \leq \phi$ if and only if there is unique a $T \in \pi(\mathfrak{A})'$ with $0 \leq T \leq 1$, s.t. $\psi(a) = \langle \pi(a)Te, e \rangle$ for all $a \in \mathfrak{A}$.

Proof. Assume there is a $T \in \pi(\mathfrak{A})'$ with $0 \leq T \leq 1$. If $a \in \mathfrak{A}_+$, then

$$0 \leq T^{\frac{1}{2}} \pi(a) T^{\frac{1}{2}} = \pi(a)T = \pi(a)^{\frac{1}{2}} T \pi(a)^{\frac{1}{2}} \leq \pi(a)$$

Therefore, $\psi(a) \leq \phi(a)$.

Conversely, we define a sesquilinear form on $\pi(\mathfrak{A})e$,

$$f(\pi(a)e, \pi(b)e) = \psi(b^*a)$$

And since $\psi \leq \phi$, by CBS Inequality, we have

$$\begin{aligned} |f(\pi(a)e, \pi(b)e)| &= |\psi(b^*a)|^2 \\ &\leq \psi(a^*a)\psi(b^*b) \leq \phi(a^*a)\phi(b^*b) \\ &= \|\pi(a)e\|^2 \|\pi(b)e\|^2 \end{aligned}$$

Therefore, f is bounded and can be extended to $\mathcal{H} = \overline{\pi(\mathfrak{A})e}$. By **Theorem 10** in the subsection **1.5.1**, we have $T \in \mathcal{B}(\mathcal{H})$, s.t.

$$\psi(b^*a) = \langle T\pi(a)e, \pi(b)e \rangle$$

In particular, $\psi(a) = \langle T\pi(a)e, e \rangle$. For $a, b, c \in \mathfrak{A}$,

$$\begin{aligned} \langle T\pi(a)\pi(b)e, \pi(c)e \rangle &= \langle T\pi(ab)e, \pi(c)e \rangle = \psi(c^*(ab)) \\ &= \psi((a^*c)^*b) = \langle T\pi(b)e, \pi(a^*c)e \rangle \\ &= \langle \pi(a)T\pi(b)e, \pi(c)e \rangle \end{aligned}$$

And by the fact that $\mathcal{H} = \overline{\pi(\mathfrak{A})e}$, $T\pi(a) = \pi(a)T$ i.e. $T \in \pi(\mathfrak{A})'$. And the uniqueness of T can be easily checked. \square

Proposition 8. *Let $(\pi_j, \mathcal{H}_j, e_j)$ for $j = 1, 2$ be two cyclic representations of a C^* -algebra \mathfrak{A} . Then $\pi_1 \cong \pi_2$ by a unitary $U: \mathcal{H}_1 \rightarrow \mathcal{H}_2$ satisfying $Ue_1 = e_2$ if and only if for any $a \in \mathfrak{A}$*

$$\langle \pi_1(a)e_1, e_1 \rangle = \langle \pi_2(a)e_2, e_2 \rangle$$

Proof. We define for any $a \in \mathfrak{A}$

$$\begin{aligned} U: \quad \pi_1(\mathfrak{A})e_1 &\longrightarrow \pi_2(\mathfrak{A})e_2 \\ \pi_1(a)e_1 &\longmapsto \pi_2(a)e_2 \end{aligned}$$

If $\langle \pi_1(a)e_1, e_1 \rangle = \langle \pi_2(a)e_2, e_2 \rangle$,

$$\|\pi_2(a)e_2\|^2 = \langle \pi_2(a)e_2, \pi_2(a)e_2 \rangle = \langle \pi_1(a)e_1, \pi_1(a)e_1 \rangle = \|\pi_1(a)e_1\|^2$$

U is a unitary on $\pi_1(\mathfrak{A})e_1$ and U can extend to \mathcal{H}_1 . And since

$$U(\pi_1(a)\pi_1(b)e_1) = \pi_2(ab)e_2 = \pi_2(a)\pi_2(b)e_2 = \pi_2(a)U(\pi_1(b)e_1)$$

U is the unitary make $\pi_1 \cong \pi_2$. The converse is trivial. \square

Remark. Using this proposition, we can strengthen the result of the **Corollary 4** in the subsection **2.2.6**. In fact, every state on an closed ideal \mathfrak{J} of a C^* -algebra \mathfrak{A} can extend to a unique state on \mathfrak{A} . If there are two extensions, they can generate two cyclic representations. And these two representations are equivalent on the ideal, thus these two representations are equivalent. And by above proposition, these two extensions are equal.

By above two propositions, we can get the relationship between positive functionals and the generated representations.

Corollary 9. *Let \mathfrak{A} be a C^* -algebra and ψ and ϕ be two positive linear functionals on \mathfrak{A} with $\psi \leq \phi$. Then π_ψ is equivalent to a subrepresentation of π_ϕ .*

Proof. By above proposition, there is a $T \in \pi(\mathfrak{A})'$ s.t.

$$\psi(a) = \langle \pi_\phi(a) T e_\phi, e_\phi \rangle$$

Put $\mathcal{M} = \overline{\pi_\phi(\mathfrak{A}) T^{\frac{1}{2}} e_\phi}$. Clearly, \mathcal{M} reduces $\pi_\phi(\mathfrak{A})$. Since

$$\langle \pi_\psi(a) e_\psi, e_\psi \rangle = \psi(a) = \langle \pi_\phi(a) T^{\frac{1}{2}} e_\phi, T^{\frac{1}{2}} e_\phi \rangle$$

by above proposition $\pi_\psi \cong \pi_\phi|_{\mathcal{M}}$ □

Definition 10. A state $\phi \in \mathcal{S}_{\mathfrak{A}}$ on some C^* -algebra \mathfrak{A} is called pure if for any positive linear functional ψ on \mathfrak{A} with $\psi \leq \phi$, there is a scalar $\lambda \in [0, 1]$ s.t. $\psi = \lambda\phi$.

Theorem 11. Let \mathfrak{A} be a C^* -algebra and (ρ, \mathcal{H}) be a representation, the following statements are equivalent. Then ρ is irreducible if and only if ρ is equivalent to some cyclic representation generated by a pure state.

Proof. Suppose ρ is irreducible. Then for any nonzero unit vector $h \in \mathcal{H}$, $\overline{\rho(\mathfrak{A})h} = \mathcal{H}$. Define

$$\begin{aligned} \phi: \quad \mathfrak{A} &\longrightarrow \mathbb{C} \\ a &\longmapsto \langle \rho(a)h, h \rangle \end{aligned}$$

Then ϕ is clearly a state. If (π, \mathcal{K}, e) is the representation generated by ϕ , then

$$\langle \pi(a)e, e \rangle = \langle \rho(a)h, h \rangle$$

Thus, by above proposition, $\pi \cong \rho$.

If there is a positive linear functional $\psi \leq \phi$, then

$$\exists T \in \pi(\mathfrak{A})', \text{ s.t. } \psi(a) = \langle \pi(a)Te, e \rangle$$

But by above theorem, $\pi(\mathfrak{A})' = \mathbb{C}$, therefore $T = \lambda \in \mathbb{C}$, $\psi = \lambda\phi$.

Conversely, if $(\rho, \mathcal{H}, e) \cong (\pi_\phi, \mathcal{H}_\phi, e_\phi)$ for some pure state ϕ ,

$$\phi(a) = \langle \rho(a)e, e \rangle = \langle \pi_\phi e_\phi, e_\phi \rangle$$

Then if a projection $P \in \rho(\mathfrak{A})'$, then

$$\psi(a) = \langle \rho(a)Pe, Pe \rangle = \langle \rho(a)Pe, e \rangle$$

is a positive linear functional and $\psi \leq \phi$, therefore $\psi = \lambda\phi$ for some $\lambda \in [0, 1]$, i.e.

$$\langle \rho(a)Pe, e \rangle = \langle Pe, \rho(a^*)e \rangle = \lambda \langle \rho(a)e, e \rangle = \langle \lambda e, \rho(a^*)e \rangle$$

Since $\mathcal{H} = \overline{\rho(\mathfrak{A})e}$, $P = \lambda$, thus $P = 1$ or 0 . Therefore ρ has no proper invariant space. □

Proposition 12. *Let \mathfrak{A} be a C^* -algebra and $\phi \in \mathcal{S}_{\mathfrak{A}}$. Then ϕ is pure if and only if ϕ is the extreme point of $\mathcal{S}_{\mathfrak{A}}$.*

Remark. This proposition can be obtained by definition. In fact, we usually use the extreme point to define the pure state.

We have similar results for pure states as general states. Firstly, we need a lemma.

Proposition 13. *Let \mathfrak{A} be a C^* -algebra and \mathfrak{I} be a closed ideal of \mathfrak{A} .*

- 1) *If (ρ, \mathcal{H}) is an irreducible representation of \mathfrak{I} , then it can extend to a unique irreducible representation $(\tilde{\rho}, \mathcal{H})$ on \mathfrak{A} .*
- 2) *If (ρ, \mathcal{H}) is an irreducible representation of \mathfrak{A} , then $(\rho|_{\mathfrak{I}}, \mathcal{H})$ is an irreducible representation of \mathfrak{I} .*

Proof. For 1), since ρ is irreducible, ρ is non-degenerate. By the **Theorem 3** in this subsection, there is a unique extension $\tilde{\rho}$ of ρ of \mathfrak{A} . And clearly,

$$\mathbb{C} \subset \tilde{\rho}(\mathfrak{A})' \subset \rho(\mathfrak{I})' = \mathbb{C}$$

$\tilde{\rho}(\mathfrak{A})' = \mathbb{C}$, i.e. $\tilde{\rho}$ is irreducible.

For 2), ρ can be generated by a pure state ϕ on \mathfrak{A} . Then $\phi|_{\mathfrak{I}}$ is a pure state on \mathfrak{I} since if $\phi|_{\mathfrak{I}} = s\psi_1 + (1-s)\psi_2$ for some $\psi_1, \psi_2 \in \mathcal{S}_{\mathfrak{I}}$ and $s \in (0, 1)$, then ψ_1, ψ_2 have extension ψ'_1, ψ'_2 on \mathfrak{A} , $\phi = s\psi'_1 + (1-s)\psi'_2$ by the uniqueness of extensions, which is a contradiction. Therefore, $\phi|_{\mathfrak{I}}$ can generate an irreducible representation κ on \mathfrak{I} and $\kappa \cong \rho|_{\mathfrak{I}}$. \square

Proposition 14. *Let \mathfrak{B} be a C^* -subalgebra of \mathfrak{A} . Then every pure state on \mathfrak{B} can extend to a pure state on \mathfrak{A} .*

Proof. Define $T: \mathcal{S}_{\mathfrak{A}} \rightarrow \mathcal{S}_{\mathfrak{B}}$ by $\phi \rightarrow \phi|_{\mathfrak{B}}$. Then T is surjective, affine and wk^* -continuous. Therefore, T maps extreme points to extreme points. Thus the proposition holds. \square

Corollary 15. *Let \mathfrak{A} be a C^* -algebra and $a \in \mathfrak{A}$. Then*

$$\|a\|^2 = \sup \{ \phi(a^*a) : \phi \text{ is a pure state on } \mathfrak{A} \}$$

Proof. If $\mathfrak{A} = C(X)$ for some compact space X , this corollary is true. And then by above proposition, any pure state on $C^*(a^*a)$ can extend to a pure state on \mathfrak{A} . Therefore, this corollary holds. \square

Theorem 16. *If \mathfrak{A} is a C^* -algebra, then there is a family of irreducible representations $\{\pi_i\}_{i \in I}$ s.t. $\oplus_{i \in I} \pi_i$ is a faithful representation of \mathfrak{A} .*

Proof. Let $\{\pi_i\}_{i \in I}$ be generated by all pure state. Then by above corollary, $\oplus_{i \in I} \pi_i$ is faithful. \square

3.3 Operators on Hilbert Spaces

We have known that each C^* -algebra can be isometrically imbedded a $\mathcal{B}(\mathcal{H})$ for some Hilbert space \mathcal{H} . Therefore, we just research operator algebras other than C^* -algebras. Firstly, we need some general theories of operators on a Hilbert space.

3.3.1 Spectrums for Operators

By using the fact that $T \rightarrow T^*$ in $\mathcal{B}(\mathcal{H})$ is bijective and isometrically antilinear, we have the following proposition.

Proposition 1. *Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$. Then*

$$\sigma(T^*) = \sigma(T)^*$$

For self-adjoint operators in $\mathcal{B}(\mathcal{H})$, we have some extral properties related to the inner product.

Proposition 2. *Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$. Then T is self-adjoint if and only if $\langle Th, h \rangle \in \mathbb{R}$ for any $h \in \mathcal{H}$.*

Proof. If $\langle Th, h \rangle \in \mathbb{R}$ for any $h \in \mathcal{H}$, then for any $h, g \in \mathcal{H}$ and $\alpha \in \mathbb{C}$,

$$\langle T(h + \alpha g), h + \alpha g \rangle = \langle Th, h \rangle + \bar{\alpha} \langle Th, g \rangle + \alpha \langle Tg, h \rangle + |\alpha|^2 \langle Tg, g \rangle \in \mathbb{R}$$

Then by taking its conjugate,

$$\begin{aligned} \bar{\alpha} \langle Th, g \rangle + \alpha \langle Tg, h \rangle &= \bar{\alpha} \langle h, Tg \rangle + \alpha \langle g, Th \rangle \\ &= \bar{\alpha} \langle T^*h, g \rangle + \alpha \langle T^*g, h \rangle \end{aligned}$$

Taking $\alpha = 1$ and i , we have

$$\begin{aligned} \langle Th, g \rangle + \langle Tg, h \rangle &= \langle T^*h, g \rangle + \langle T^*g, h \rangle \\ i \langle Th, g \rangle - i \langle Tg, h \rangle &= -i \langle T^*h, g \rangle + i \langle T^*g, h \rangle \end{aligned}$$

Therefore, $\langle Th, g \rangle = \langle T^*g, h \rangle$ for any $h, g \in \mathcal{H}$.

The converse is clearly true by definition. \square

Proposition 3. *Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$ be self-adjoint. Then*

$$\sigma(T) \subset \{ \langle Th, h \rangle : \|h\| = 1 \}$$

Proof. If $\lambda \notin \{ \langle Th, h \rangle : \|h\| = 1 \}$, then there is a $\varepsilon > 0$, s.t.

$$\varepsilon < |\lambda - \langle Th, h \rangle| \leq \|Th\|, \quad \forall h \in \mathcal{H}, \quad \|h\| = 1$$

Therefore, by the **Lemma 3** in the subsection 1.4.4,

$$\ker(T - \lambda) = \{0\} \text{ \& \; } \text{ran}(T - \lambda) \text{ is closed}$$

Since T is self-adjoint, $\bar{\lambda} \notin \sigma(T)$, so $\ker(\bar{\lambda} - T) = \{0\}$. Thus

$$\operatorname{ran}(T - \lambda)^\perp = \ker(T - \bar{\lambda}) = \{0\}$$

i.e. $\operatorname{ran}(T - \lambda) = \mathcal{H}$. Then by the Inverse Mapping Theorem, $(T - \lambda)^{-1} \in \mathcal{B}(\mathcal{H})$, i.e. $\lambda \notin \sigma(T)$. \square

Corollary 4. *Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$ be self-adjoint.*

- 1) $\|T\| = \sup \{ \langle Th, h \rangle : \|h\| = 1 \}$.
- 2) If $\langle Th, h \rangle = 0$ for any $h \in \mathcal{H}$, then $T = 0$.

Proof. For 1),

$$\|T\| = r(T) \leq \sup \{ \langle Th, h \rangle : \|h\| = 1 \} \leq \|T\|$$

2) is the direct result of 1). \square

Remark. In fact, for any $T \in \mathcal{B}(\mathcal{H})$ with $\langle Th, h \rangle = 0$ for any $h \in \mathcal{H}$, then $T = 0$ since any T can be the linear combination of two self-adjoint operators.

By above proposition, we see the spectrum of operators can be described by using another method. Firstly, we can generalize the consequence in **Proposition 3**.

Proposition 5. *Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$. Then let*

$$\chi(\lambda) = \inf \{ \|(T - \lambda)h\| : \|h\| = 1 \}$$

$$\sigma_l(T) = \{ \lambda \in \mathbb{C} : \chi(\lambda) = 0 \}$$

Proof. By above proposition, we have seen if $\chi(\lambda) > 0$,

$$\ker(T - \lambda) = \{0\} \text{ \& \; } \operatorname{ran}(T - \lambda) \text{ is closed}$$

Then let $\mathcal{K} = \operatorname{ran}(T - \lambda)$, then $T - \lambda : \mathcal{H} \rightarrow \mathcal{K}$ is a bijection and $(T - \lambda)^{-1} : \mathcal{K} \rightarrow \mathcal{H}$ is bounded. Let

$$B = (T - \lambda)^{-1}P : \mathcal{K} \rightarrow \mathcal{H}$$

where P is a projection to \mathcal{K} . Then $B \in \mathcal{B}(\mathcal{H})$ and $B(T - \lambda) = 1$. Therefore, $\sigma_l(T) \subset \{ \lambda \in \mathbb{C} : \chi(\lambda) = 0 \}$.

Conversely, if there is a $B \in \mathcal{B}(\mathcal{H})$ s.t. $B(T - \lambda) = 1$, then for $h \in \mathcal{H}$ with $\|h\| = 1$

$$1 = \|h\| = \|B(T - \lambda)h\| \leq \|B\| \|(T - \lambda)h\|$$

Therefore, $\|(T - \lambda)h\| \geq \|B\|^{-1}$, i.e. $\chi(\lambda) > 0$. \square

Corollary 6. *If $N \in \mathcal{B}(\mathcal{H})$ is a normal operator, then*

$$\sigma_l(N) = \sigma_r(N) = \sigma(N)$$

In particular, if a normal operator is left (or right) invertible, then it is invertible.

Proof. If N is normal, then $N - \lambda$ is normal. Therefore,

$$\|(N - \lambda)h\| = \|(N^* - \bar{\lambda})h\|, \quad \forall h \in \mathcal{H}$$

By above theorem, $\sigma_l(N^*) = \overline{\sigma_l(N)}$. Thus,

$$\sigma_r(N) = \overline{\sigma_l(N^*)} = \sigma_l(N) \quad \square$$

Definition 7. Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$.

1) The point spectrum of T is defined as

$$\sigma_p(T) = \{ \lambda \in \mathbb{C} : \ker(T - \lambda) \neq \{0\} \}$$

2) The approximate point spectrum of T is defined as

$$\sigma_{ap}(T) = \{ \lambda \in \mathbb{C} : \exists \text{ units } x_n \subset \mathcal{B}(\mathcal{H}) \text{ s.t. } \|(T - \lambda)x_n\| \rightarrow 0 \}$$

Proposition 8. *Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$.*

1. $\lambda \notin \sigma_{ap}(T)$ if and only if $\ker(T - \lambda) = \{0\}$ and $\text{ran}(T - \lambda)$ is closed.
2. $\partial\sigma(T) \subset \sigma_{ap}(T)$.
3. $\partial\sigma(T) \subset \sigma_l(T) \cap \sigma_r(T) = \sigma_{ap}(T) \cap \sigma_{ap}(T^*)$.

Proof. For 1), if $\lambda \notin \sigma_{ap}(T)$, then there is a constant $C > 0$ s.t.

$$\|(T - \lambda)x\| \geq C \|x\|, \quad \forall x \in \mathcal{H}$$

Thus by above proposition, $\ker(T - \lambda) = \{0\}$ and $\text{ran}(T - \lambda)$ is closed. Conversely, the condition that means $(T - \lambda)$ is left invertible, thus

$$\|(T - \lambda)x\| \geq \|B\|^{-1} \|x\|, \quad \forall x \in \mathcal{H}, \quad \|x\| = 1$$

where B is the left inverse of T . Therefore, $\lambda \notin \sigma_{ap}(T)$.

For 2), let $\lambda \in \partial\sigma(T)$ and $\{\lambda_n\} \subset \rho(T)$ s.t. $\lambda_n \rightarrow \lambda$. Then we know $\|(T - \lambda_n)^{-1}\| \rightarrow \infty$. Otherwise, if $\|(T - \lambda_n)^{-1}\| < M$ and $|\lambda_n - \lambda| < M^{-1}$ for sufficiently large n ,

$$\|(T - \lambda_n) - (T - \lambda)\| < \|(T - \lambda_n)^{-1}\|^{-1}$$

by the perturbation of inverse, $T - \lambda$ is also invertible. Let $\{x_n\}$ by any sequence of unit vectors in \mathcal{H} . Define

$$y_n = \frac{(T - \lambda_n)^{-1}x_n}{\|(T - \lambda_n)^{-1}x_n\|} \quad (3.1)$$

Then $\|(T - \lambda)y_n\| \rightarrow 0$, i.e. $\lambda \in \sigma_{ap}(T)$.

For 3), by 1) we have known $\sigma_l(T) = \sigma_{ap}(T)$ and by taking the complex conjugate, $\sigma_r(T^*) = \sigma_{ap}(T^*)$. Therefore, 3) holds. \square

Now, we can prove the inverse of the **Corollary 6** in the subsection **2.2.3**.

Theorem 9. *Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$. If for any $h \in \mathcal{H}$,*

$$\langle Th, h \rangle \geq 0$$

then T is positive.

Proof. Firstly, T is self-adjoint by above proposition. And for $\lambda < 0$ and $\|h\| = 1$

$$\begin{aligned} \|(T - \lambda)h\|^2 &= \|Th\|^2 - 2\lambda\langle Th, h \rangle + \lambda^2\|h\|^2 \\ &\geq -2\lambda\langle Th, h \rangle + \lambda^2\|h\|^2 \\ &\geq \lambda^2\|h\|^2 \end{aligned}$$

Therefore, $\lambda \notin \sigma_l(T)$. Since T is self-adjoint, $\lambda \notin \sigma_r(T)$. Thus $\lambda \notin \sigma(T)$, i.e. $\sigma(T) \subset [0, \infty)$. \square

3.3.2 Polar Decomposition

For a $z \in \mathbb{C}$, it can be decomposed as $z = |z|e^{i\theta}$. Similarly, a operator in $\mathcal{B}(\mathcal{H})$ can also be decomposed as a complex number does. Firstly, we need to find the element playing same role as $e^{i\theta}$.

Definition 1. Let \mathcal{H} be a Hilbert space and $U \in \mathcal{B}(\mathcal{H})$ and \mathcal{M} be a closed subspace of \mathcal{H} . If $U|_{\mathcal{M}}$ is an isometry and $U|_{\mathcal{M}^\perp} = 0$, then U is called a partial isometry. \mathcal{M} is called the initial space and $\mathcal{N} = U(\mathcal{M})$ is called the final space.

Proposition 2. *If U is a partial isometry, then U^*U is the projection on the initial space \mathcal{M} and UU^* is the projection on the final space \mathcal{N} . Moreover, U^* is the partial isometry with the initial space \mathcal{N} and final space \mathcal{M} .*

Proof. For any $g, h \in \mathcal{H}$, then $g = g_{\parallel} + g_{\perp}$ and $h = h_{\parallel} + h_{\perp}$, where $g_{\parallel}, h_{\parallel} \in \mathcal{M}$ and $g_{\perp}, h_{\perp} \in \mathcal{M}^\perp$

$$\begin{aligned} \langle U^*U(g_{\parallel} + g_{\perp}), h_{\parallel} + h_{\perp} \rangle &= \langle Ug_{\parallel}, Uh_{\parallel} \rangle = \langle g_{\parallel}, h_{\parallel} \rangle \\ &= \langle P_{\mathcal{M}}(g_{\parallel} + g_{\perp}), h_{\parallel} + h_{\perp} \rangle \end{aligned}$$

Then $U^*U = P_{\mathcal{M}}$. $P_{\mathcal{M}^\perp} = 1 - U^*U$ implies $U(1 - U^*U) = 0$, i.e. $U = UU^*U$ and $\ker U^* = (\text{ran } U)^\perp = \mathcal{N}^\perp$, therefore, UU^* is the projection on \mathcal{N} . For $x \in \mathcal{N}$.

$$\langle U^*x, U^*x \rangle = \langle UU^*x, x \rangle = \langle x, x \rangle \quad (3.2)$$

And $\text{ran } U^* = (\ker U)^\perp = \mathcal{M}$, thus U^* is a partial isometry with the initial space \mathcal{N} and the final space \mathcal{M} . \square

Corollary 3. *If $U \in \mathcal{B}(\mathcal{H})$, then the following statements are equivalent.*

- 1) U is a partial isometry.
- 2) U^*U is a projection.
- 3) $U = UU^*U$.
- 4) U^* is a partial isometry.
- 5) UU^* is a projection.
- 6) $U^* = U^*UU^*$.

Theorem 4 (Polar Decomposition). *Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$. Then there is a partial isometry $U \in \mathcal{B}(\mathcal{H})$ with the initial space $(\ker T)^\perp$ and the final space $\overline{\text{ran } T}$, s.t. $T = U|T|$. Moreover, if $T = WS$, where $S \geq 0$ and W is a partial isometry with $\ker W = \ker S$, then $S = |T|$ and $W = U$.*

Proof. Define

$$\begin{aligned} U: \quad \text{ran } |T| &\longrightarrow \text{ran } T \\ |T| h &\longmapsto Th \end{aligned}$$

Since

$$\|Th\|^2 = \langle Th, Th \rangle = \langle T^*Th, h \rangle = \langle |T| h, |T| h \rangle = \||T| h\|^2$$

U is a isometry. And thus U can extend to $\overline{\text{ran } |T|}$ and onto $\overline{\text{ran } T}$. Then extending U to \mathcal{H} by setting $U|_{\overline{\text{ran } |T|}^\perp} = 0$. U is a partial isometry with the initial space $\ker |T|^\perp$ and the final space $\overline{\text{ran } T}$. Clearly, $\ker |T| = \ker T$ and $T = U|T|$.

If $T = WS$, then $T^*T = SW^*WS = SP_{\ker S^\perp}S = S^2$. Therefore, $S = |T|$. And $W|T| = U|T|$, thus W and U has same initial space, then $W = U$. \square

3.3.3 Topologies on $\mathcal{B}(\mathcal{H})$

In $\mathcal{B}(\mathcal{H})$, the norm topology may be too stronger, which means norm closed subalgebras in $\mathcal{B}(\mathcal{H})$ do not contain "enough" elements. Therefore, we need some weaker topologies

Definition 1. Let \mathcal{H} be a Hilbert space.

- 1) The strong operator topology on $\mathcal{B}(\mathcal{H})$, denoted by SOT , is generated by a family of seminorms $\{p_x\}_{x \in \mathcal{H}}$, where $p_x(T) = \|Tx\|$.
- 2) The weak operator topology on $\mathcal{B}(\mathcal{H})$, denoted by WOT , is generated by a family of seminorms $\{p_{gh}\}_{g,h \in \mathcal{H}}$, where $p_{gh}(T) = |\langle Tg, h \rangle|$.

Remark. Firstly, by definition and the **Corollary 4** in the subsection **2.3.1**, we can see

$$\bigcap_{x \in \mathcal{H}} \ker p_x = \{0\} \ \& \ \bigcap_{g, h \in \mathcal{H}} \ker p_{gh} = \{0\}$$

Therefore, both *SOT* and *WOT* can let $\mathcal{B}(\mathcal{H})$ be a Hausdorff locally convex t.v.s.. And by the CBS Inequality, *SOT* is stronger than *WOT*. Moreover,

1) the subbasis of the *WOT* at T_0 is like

$$U_{gh\varepsilon}(T_0) = \{T \in \mathcal{B}(\mathcal{H}) : |\langle (T - T_0)g, h \rangle| < \varepsilon\}$$

Therefore, a net $\{T_\alpha\}$ in $\mathcal{B}(\mathcal{H})$ converges in *WOT* to T_0 if and only if $\langle T_\alpha g, h \rangle \rightarrow \langle T_0 g, h \rangle$ for all $g, h \in \mathcal{H}$.

2) the subbasis of the *SOT* at T_0 is like

$$V_{x\varepsilon}(T_0) = \{T \in \mathcal{B}(\mathcal{H}) : \|(T - T_0)x\| < \varepsilon\}$$

Therefore, a net $\{T_\alpha\}$ in $\mathcal{B}(\mathcal{H})$ converges in *SOT* to T_0 if and only if $\|T_\alpha x\| \rightarrow \|T_0 x\|$ for all $x \in \mathcal{H}$.

Comparing with the norm topology, the *SOT* is related to the pointwise convergence, where the norm topology is related to the uniform convergence.

Proposition 2. *Let \mathcal{H} be a Hilbert space and L is a linear functional on $\mathcal{B}(\mathcal{H})$. Then the following statements are equivalent.*

- 1) L is *SOT*-continuous.
- 2) L is *WOT*-continuous.
- 3) There are vectors $f_1, \dots, f_n, g_1, \dots, g_n \in \mathcal{H}$ s.t.

$$L(T) = \sum_{i=1}^n \langle Tg_i, f_i \rangle, \quad \forall T \in \mathcal{B}(\mathcal{H})$$

Proof. Clearly, 3) \Rightarrow 2) \Rightarrow 1), thus assume 1), then by the **Theorem 4** in the subsection **1.2.4**, there are $g_1, \dots, g_n \in \mathcal{H}$ and $\alpha_1, \dots, \alpha_n > 0$ s.t.

$$\|L(T)\| \leq \sum_{i=1}^n \alpha_i \|Tg_i\|$$

Replacing g_i by $\alpha_i g_i$ and g_i by $\sqrt{n}g_i$ then

$$\|L(T)\| \leq \sum_{i=1}^n \|Tg_i\| \leq \sqrt{n} \left(\sum_{i=1}^n \|Tg_i\|^2 \right)^{\frac{1}{2}} = \left(\sum_{i=1}^n \|Tg_i\|^2 \right)^{\frac{1}{2}}$$

Let $\mathcal{K} = \overline{\{\oplus_{i=1}^n Tg_i : T \in \mathcal{B}(\mathcal{H})\}}$ be a closed subspace of $\mathcal{H}^{(n)}$. Define a linear functional F on \mathcal{K} by $F(\oplus_{i=1}^n Tg_i) = L(T)$, since

$$|F(\oplus_{i=1}^n Tg_i)| = \|L(T)\| \leq \left(\sum_{i=1}^n \|Tg_i\|^2\right)^{\frac{1}{2}}$$

Therefore, F can extend to $\mathcal{H}^{(n)}$ with the coincided inner product and norm. By the Riesz Theorem, there is a $f_1, \dots, f_n \in \mathcal{H}$ s.t.

$$F(\oplus_{i=1}^n h_i) = \langle \oplus_{i=1}^n h_i, \oplus_{i=1}^n f_i \rangle, \quad \forall \oplus_{i=1}^n h_i \in \mathcal{K}^{(n)}$$

In particular, restricting F to \mathcal{K}

$$L(T) = F(\oplus_{i=1}^n Tg_i) = \langle \oplus_{i=1}^n Tg_i, \oplus_{i=1}^n f_i \rangle = \sum_{i=1}^n \langle Tg_i, f_i \rangle \quad \square$$

Proposition 3. *If S is a convex subset of $\mathcal{B}(\mathcal{H})$, then*

$$\overline{S}^{WOT} = \overline{S}^{SOT}$$

Proof. Clearly,

$$\mathcal{T}_{WOT} \subset \mathcal{T}_{SOT}$$

Therefore,

$$\overline{S}^{SOT} \subset \overline{S}^{WOT}$$

Since for linear functionals, SOT -continuity implies WOT -continuity,

$$\overline{(S, WOT)}^{wk} \subset \overline{(S, SOT)}^{wk}$$

And we have known $\overline{(S, WOT)}^{wk} = \overline{(S, WOT)} = \overline{S}^{WOT}$ and $\overline{(S, SOT)}^{wk} = \overline{(S, SOT)} = \overline{S}^{SOT}$, therefore

$$\overline{S}^{SOT} \subset \overline{S}^{WOT} \subset \overline{S}^{SOT} \quad \square$$

WOT also have a similar property to wk -topology, the Alaoglu's Theorem.

Theorem 4. *The closed unit ball in $\mathcal{B}(\mathcal{H})$ is WOT -compact.*

Proof. Let X_h be closed unit ball in \mathcal{H} with the wk -topology and indexed with $h \in \mathcal{H}$ and B be the closed unit ball in \mathcal{H} with the WOT . Define

$$\begin{aligned} \tau: \quad B &\longrightarrow \prod \{X_h : \|h\| \leq 1\} \\ T &\longmapsto \tau(T) \end{aligned}$$

where $\tau(T)_h = Th$. Then using similar method, we can prove B is compact. \square

Chapter 4

Normal Operators and Abelian Von Neumann Algebras

For a normal operator N , it can generate an abelian norm closed C^* -subalgebra of $\mathcal{B}(\mathcal{H})$, which is isometrically $*$ -isomorphic to $C(\sigma(N))$. But the norm closed C^* -subalgebra may be too "small" to get more interesting results. Therefore, we want to see the WOT -closed C^* -subalgebra $W^*(N)$ generated by a normal operator N . This extension can be naturally done by following processes. Firstly, we will extend $C(X)$ to $B(X, \Omega)$, where $B(X, \Omega)$ denotes all Ω -measurable functions on X , and find the map from $B(X, \Omega)$ to $\mathcal{B}(\mathcal{H})$. Then using these method, we can extend the Continuous Functional Calculus to the Borel Functional Calculus. In order to illuminate the structure of $W^*(N)$, we need more properties of WOT -closed C^* -subalgebras, i.e. von Neumann algebras, and abelian von Neumann algebras. After that, we can build the isomorphism from $W^*(N)$ to $L^\infty(\sigma(N))$. Finally, by above methods, we can get the multiplicity theory of normal operators.

4.1 Spectral Theorem

4.1.1 Spectral Measures

Firstly, we want to construct the map from $B(X, \Omega)$ to $\mathcal{B}(\mathcal{H})$. Thus we need to define the operator-valued measure.

Definition 1. If X is a set and Ω is a σ -algebra of subsets of X , and \mathcal{H} is a Hilbert space, then a spectral measure for (X, Ω, \mathcal{H}) is a function $E: \Omega \rightarrow \mathcal{B}(\mathcal{H})$ satisfying the following properties.

- 1) $E(\Delta)$ is a projection for any $\Delta \in \Omega$.
- 2) $E(\emptyset) = 0$ and $E(X) = 1$.

3) $E(\Delta_1 \cap \Delta_2) = E(\Delta_1)E(\Delta_2)$ for any $\Delta_1, \Delta_2 \in \Omega$.

4) If $\{\Delta_n\}_{n=1}^\infty$ is a sequence of pairwise disjoint sets in Ω ,

$$E\left(\bigcup_{n=1}^\infty \Delta_n\right) = \sum_{n=1}^\infty E(\Delta_n)$$

Remark. In 4), this convergence is about *SOT*, and since $E(\Delta)$ is a projection, it is also about *WOT*.

Then we can see the relation between the spectral measure and the complex measure.

Proposition 2. *If E is a spectral measure for (X, Ω, \mathcal{H}) and $g, h \in \mathcal{H}$, then*

$$E_{gh}(\Delta) = \langle E(\Delta)g, h \rangle$$

is a complex measure of (X, Ω) . Moreover, $\|E_{gh}\| \leq \|g\| \|h\|$.

Proof. $\mu = E_{gh}(\Delta)$ is a complex measure by definition.

$$|\mu|(\Delta) = |\langle E(\Delta)g, h \rangle| \leq \|E(\Delta)g\| \|h\| \leq \|g\| \|h\| \quad \square$$

Using the spectral measure, we can define the operator-valued integral.

Theorem 3. *If E is a spectral measure for (X, Ω, \mathcal{H}) and $\phi: X \rightarrow \mathbb{C}$ is a bounded Ω -measurable function, then there is a unique operator $T \in \mathcal{B}(\mathcal{H})$ s.t. for any $\varepsilon > 0$, there is a Ω -partition $\{\Delta_i\}_{i=1}^n$ of X with*

$$\sup \{|\phi(x) - \phi(y)| : x, y \in \Delta_k\}$$

for $1 \leq k \leq n$ s.t. for any $x_k \in \Delta_k$

$$\left\| T - \sum_{k=1}^n \phi(x_k) E(\Delta_k) \right\| < \varepsilon$$

Proof. Define

$$\begin{aligned} B: \quad \mathcal{H} \times \mathcal{H} &\longrightarrow \mathbb{C} \\ (g, h) &\longmapsto \int \phi dE_{gh} \end{aligned}$$

Since $\|B(g, h)\| \leq \|\phi\|_\infty \|g\| \|h\|$, B is a bounded sesquilinear form. Then by the Riesz Theorem, there is a $B \in \mathcal{B}(\mathcal{H})$ s.t. $B(g, h) = \langle Tg, h \rangle$ for all $g, h \in \mathcal{H}$ and $\|T\| \leq \|\phi\|_\infty$. Then for and $g, h \in \mathcal{H}$ and $x_k \in \Delta_k$ with the giving partition

$$\begin{aligned} \left| \langle Tg, h \rangle - \sum_{k=1}^n \phi(x_k) \langle E(\Delta_k)g, h \rangle \right| &= \left| \sum_{k=1}^n \int_{\Delta_k} (\phi(x) - \phi(x_k)) dE_{gh}(x) \right| \\ &\leq \sum_{k=1}^n \int_{\Delta_k} |\phi(x) - \phi(x_k)| d|E_{gh}|(x) \\ &\leq \varepsilon \int d|E_{gh}|(x) \leq \varepsilon \|g\| \|h\| \end{aligned}$$

Therefore, by the Riezs Theorem,

$$\left\| T - \sum_{k=1}^n \phi(k) E(\Delta_k) \right\| < \varepsilon \quad \square$$

Remark. We define $T = \int \phi dE$, and $\langle Tg, h \rangle = \int \phi dE_{gh}$.

We have one more property on $B(X)$ for a compact space X . Let $B(X)$ denote all bounded Borel measurable functions on X and $M(X)$ be all Borel measures.

Lemma 4. *If X is a compact space and $\phi \in B(X)$, then there is a net $\{u_i\} \subset C(X)$ with $\|u_i\|_\infty \leq \|\phi\|$ s.t. $\int u_i d\mu \rightarrow \int \phi d\mu$ for any μ in $M(X)$.*

Proof. $C(X)^* = M(X)$ and $C(X) \subset C(X)^{**} = M(X)^*$. Therefore, the unit ball in $C(X)$ is wk -dense in the unit ball in $M(X)^*$ (by using Hahn-Banach Theorem). Identifying $B(X)$ be the subspace of $M(X)^*$, we have this lemma. \square

Remark. In fact, $C(X)$ is norm closed in $C(X)^{**} = M(X)^*$ like that the C^* -algebra is norm closed in $\mathcal{B}(\mathcal{H})$. But by above lemma, we see in $M(X)^*$

$$\overline{C(X)}^{wk} = B(X)$$

Thus we may apply this to extending C^* -algebras.

$B(X, \Omega)$ with the supremum norm and complex conjugate can become a C^* -algebra. Then by above theorem, we can provide the map from $B(X, \Omega)$ to operator algebras.

Proposition 5. *If E is a spectral measure for (X, Ω, \mathcal{H}) , then*

$$\begin{aligned} \rho: \quad B(X, \Omega) &\longrightarrow \mathcal{B}(\mathcal{H}) \\ \phi &\longmapsto \int \phi dE \end{aligned}$$

is a representation of $B(X, \Omega)$.

Proof. Firstly, we need check that ρ preserves the involution. For $g, h \in \mathcal{H}$,

$$E_{gh}(\Delta) = \langle E(\Delta)g, h \rangle = \langle g, E(\Delta)h \rangle = \overline{\langle E(\Delta)h, g \rangle} = \overline{E_{hg}(\Delta)}$$

Therefore, for any $g, h \in \mathcal{H}$

$$\begin{aligned} \langle (\int \phi dE)^* g, h \rangle &= \langle g, (\int \phi dE) h \rangle \\ &= \overline{\langle (\int \phi dE) h, g \rangle} = \overline{\int \phi dE_{hg}} \\ &= \int \phi^* d\overline{E_{hg}} = \int \phi^* dE_{gh} \\ &= \langle (\int \phi^* dE) g, h \rangle \end{aligned}$$

Also, ρ must be multiplicative. Fix $\phi, \psi \in B(X)$ and $\varepsilon > 0$, let $\{\Delta_i\}_{i=1}^n$ be the Ω -partition of X s.t. the oscillations of $\phi, \psi, \phi\psi$ on each Δ_i are less than ε . By the fact $E(\Delta_i)E(\Delta_j) = 0$ for $i \neq j$,

$$\begin{aligned}
& \left\| \int \phi\psi dE - \left(\int \phi dE \right) \left(\int \psi dE \right) \right\| \\
& \leq \varepsilon + \left\| \sum_{k=1}^n \phi(x_k)\psi(x_k)E(\Delta_k) - \left(\sum_{i=1}^n \phi(x_i)E(\Delta_i) \right) \left(\sum_{j=1}^n \psi(x_j)E(\Delta_j) \right) \right\| \\
& \quad + \left\| \left(\sum_{i=1}^n \phi(x_i)E(\Delta_i) \right) \left(\sum_{j=1}^n \psi(x_j)E(\Delta_j) \right) - \left(\int \phi dE \right) \left(\int \psi dE \right) \right\| \\
& \leq \varepsilon + \left\| \left(\sum_{i=1}^n \phi(x_i)E(\Delta_i) \right) \left(\sum_{j=1}^n \psi(x_j)E(\Delta_j) - \int \psi dE \right) \right\| \\
& \quad + \left\| \left(\sum_{i=1}^n \phi(x_i)E(\Delta_i) - \int \phi dE \right) \left(\sum_{j=1}^n \psi(x_j)E(\Delta_j) \right) \right\| \\
& \leq \varepsilon(1 + \|\phi\|_\infty + \|\psi\|_\infty)
\end{aligned}$$

Then ρ is indeed multiplicative. \square

We have already there is a map from $C(X)$ to $\mathcal{B}(\mathcal{H})$, then by combining this proposition and above lemma, if the map is wk -continuous on $C(X)$, we can extend it from $M(X)$ to $\mathcal{B}(\mathcal{H})$.

Theorem 6. *If X is a compact space and $\rho: C(X) \rightarrow \mathcal{B}(\mathcal{H})$ is a representation, then there is a unique spectral measure E defined on the Borel sets of X s.t. ρ can be expressed as*

$$\rho(\phi) = \int \phi dE \quad \forall \phi \in C(X)$$

Moreover, $T \in \mathcal{B}(\mathcal{H})$ commutes with $\rho(\phi)$ for any $\phi \in C(X)$ if and only if T commutes with $E(\Delta)$ for any Borel set Δ .

Proof. Constructing $\tilde{\rho}$ from ρ :

For any $g, h \in \mathcal{H}$, the map $u \rightarrow \langle \rho(u)g, h \rangle$ from $C(X)$ to \mathbb{C} is bounded with $\|g\| \|h\|$. Therefore, by Riesz Theorem, there exists a unique $\mu_{gh} \in M(X)$ with $\|\mu_{gh}\| \leq \|g\| \|h\|$

$$\langle \rho(u)g, h \rangle = \int u d\mu_{gh}, \quad \forall u \in C(X)$$

For any $\phi \in B(X)$, $(g, h) \rightarrow \int \phi d\mu_{gh}$ is a bounded sesquilinear on $\mathcal{H} \times \mathcal{H}$. Therefore, there exists a unique $T \in \mathcal{B}(\mathcal{H})$, s.t.

$$\int \phi d\mu_{gh} = \langle Tg, h \rangle$$

Then we define $\tilde{\rho}(\phi) = T$ and it is well-defined. Clearly, $\tilde{\rho}(u) = \rho(u)$ for any u in $C(X)$.

Check: $\tilde{\rho}: B(X) \rightarrow \mathcal{B}(\mathcal{H})$ is a representation

By the definition, $\tilde{\rho}$ is clearly linear and preserves the involution, thus we just need to prove the multiplication. Let ϕ and ψ in $B(X)$. There is a sequence $\{u_i\} \subset C(X)$ s.t. $u_i \rightarrow \phi$ in wk -topology. Then since $\int u_i \psi d\mu \rightarrow \int \phi \psi$ for any $\mu \in M(X)$, $\rho(u_i \psi) \rightarrow \rho(\phi \psi)$ in WOT . For any $\psi \in C(X)$,

$$\tilde{\rho}(\phi \psi) = \lim \tilde{\rho}(u_i \psi) = \lim \rho(u_i) \rho(\psi) = \tilde{\rho}(\phi) \rho(\psi)$$

Therefore, if $\phi, \psi \in B(X)$, $\tilde{\rho}(\phi \psi) = \tilde{\rho}(\phi) \tilde{\rho}(\psi)$.

Constructing the spectral measure:

Let $E(\Delta) = \rho(\chi_\Delta(x))$ for all Borel set Δ . Clearly, $E(\Delta)$ is a projection for any Δ . For disjoint Borel sets Δ_1, Δ_2 ,

$$E(\Delta_1 \cap \Delta_2) = \tilde{\rho}(\chi_{\Delta_1 \cap \Delta_2}) = \tilde{\rho}(\chi(\Delta_1) \chi(\Delta_2)) = \tilde{\rho}(\chi(\Delta_1)) = E(\Delta_1) E(\Delta_2)$$

Let $\{\Delta_n\}_{n=1}^\infty$ be disjoint Borel sets and $\Delta = \cup_{n=1}^\infty \Delta_n$ and $\Lambda_n = \cup_{k=n+1}^\infty \Delta_k$. By the fact that E is finitely additive, for any $h \in \mathcal{H}$

$$\begin{aligned} \left\| E(\Delta)h - \sum_{k=1}^n E(\Delta_k)h \right\| &= \langle E(\Lambda_n)h, E(\Lambda_n)h \rangle \\ &= \langle E(\Lambda_n)h, h \rangle = \langle \tilde{\rho}(\Lambda_n)h, h \rangle \\ &= \int \chi_{\Lambda_n} d\mu_{hh} \\ &\leq \mu_{hh}(\Lambda_n) \rightarrow 0 \end{aligned}$$

Then by above theorem, with respect to this spectral measure E , for any $\phi \in B(X)$,

$$\tilde{\rho}(\phi) = \int \phi dE$$

In particular, for $u \in C(X)$, it is also true.

Check: The uniqueness of E . If there is another E' satisfying above conditions, then for any $g, h \in \mathcal{H}$

$$\int \phi dE'_{gh} = \int \phi dE_{gh}$$

Therefore, $E' = E$.

Finally, if T commutes with $\tilde{\rho}(\phi)$ for any $\phi \in C(X)$, we can use the WOT -convergence to get that T commutes with $\tilde{\rho}(\phi)$ for any $\phi \in B(X)$. Therefore, clearly T commutes with $E(\Delta)$ for any Borel set Δ . Conversely, if T commutes with $E(\Delta)$ for any Borel set Δ ,

$$\langle E(\Delta)Tg, h \rangle = \langle E(\Delta)g, T^*h \rangle$$

i.e. $\mu_{Tg,h} = \mu_{g,T^*h}$. Thus for any $\phi \in B(X)$,

$$\int \phi d\mu_{Tg,h} = \int \phi d\mu_{g,T^*h}, \text{ i.e. } \langle \phi Tg, h \rangle = \langle \phi g, T^*h \rangle$$

Then $\langle \tilde{\rho}(\phi)Tg, h \rangle = \langle T\tilde{\rho}(\phi)g, h \rangle$ for any $g, h \in \mathcal{H}$. \square

Remark. Thus, we extend $\rho: C(X) \rightarrow \mathcal{B}(\mathcal{H})$ to $\tilde{\rho}$ defined on $B(X)$. Moreover, by above proof we know the $*$ -homomorphism

$$\begin{aligned} \tilde{\rho}: (B(X), wk) &\longrightarrow (\mathcal{B}(\mathcal{H}), WOT) \\ \phi &\longmapsto \int \phi dE \end{aligned}$$

is continuous and therefore $\tilde{\rho}(B(X))$ is WOT -closed in $\mathcal{B}(\mathcal{H})$.

4.1.2 Spectral Theorem for Normal Operators

Now we can apply above theorem to a normal operator and extend the Continuous Functional Calculus to the Borel Functional Calculus. Moreover, it will promote us to consider the structure of the WOT -closed operator algebra generated by a normal operator.

Theorem 1. *If N is a normal operator, then there is a unique spectral measure E s.t.*

$$\begin{aligned} \rho: B(\sigma(N)) &\longrightarrow \mathcal{B}(\mathcal{H}) \\ \phi &\longmapsto \phi(N) = \int \phi dE \end{aligned}$$

and $\rho|_{C(\sigma(N))}$ is the Continuous Functional Calculus. Moreover,

- 1) $\tilde{\rho}: (B(X), wk) \rightarrow (\mathcal{B}(\mathcal{H}), WOT)$ is continuous.
- 2) if G is nonempty and relatively open subset in $\sigma(N)$, $E(G) \neq 0$.
- 3) if $T \in \mathcal{B}(\mathcal{H})$, then $TN = NT$ and $TN^* = N^*T$ if and only if T commutes with $E(\Delta)$ for any Borel set Δ .

Proof. By above theorem, this ρ is in fact the extension of the Continuous Functional Calculus. And ρ is $wk - WOT$ continuous.

For 2), it is because that $\rho|_{C(\sigma(N))}$ is an isometry. If $\chi_G > 0$, then there is a continuous function s.t. $0 < f < \chi_G$. Therefore, $0 < \rho(f) \leq \rho(\chi_G)$.

Since the condition that T commutes with N and N^* is equivalent to that T commutes with $\phi(N)$ for any $\phi \in C(\sigma(N))$, by above theorem, 3) holds. \square

Remark. Now, for a normal operator N , we have

$$N = \int_{\sigma(N)} z dE = z(N) = \int_{\Gamma} (z - N)^{-1} dz$$

There are some properties of the spectral measure of normal operator N .

Proposition 2. *Let N be a normal operator with the spectral measure E .*

- 1) $\lambda \in \sigma_p(N) \Leftrightarrow E(\{\lambda\}) \neq 0$.
- 2) If $\lambda \in \sigma_p(N)$, then $\text{ran } E(\{\lambda\}) = \ker(N - \lambda)$.
- 3) If $\lambda \in \sigma_p(N)$ and $\phi \in B(\sigma(N))$, then for any $h \in \mathcal{H}$,

$$\phi(N)h = \phi(\lambda)h$$

Proof. For 1), if $\lambda \in \sigma(N)$, then there is a nonzero vector $h \in \mathcal{H}$ s.t. $(N - \lambda)h = 0$.

$$0 = \|(N - \lambda)h\|^2 = \int_{\sigma(N)} |z - \lambda|^2 dE_{hh} = \int_{\sigma(N) \setminus \{\lambda\}} |z - \lambda|^2 dE_{hh}$$

Then since E_{hh} is a positive measure,

$$\int_{\sigma(N) \setminus \{\lambda\}} dE_{hh} = \|E(\sigma(N) \setminus \{\lambda\})h\|^2 = 0$$

And since $E(\sigma(N)) = 1$,

$$E(\{\lambda\})h = h, \text{ i.e. } E(\{\lambda\}) \neq 0$$

Conversely, if $E(\{\lambda\}) \neq 0$, then for nonzero $h \in \text{ran } E(\{\lambda\})$, $E(\{\lambda\})h = h$,

$$\begin{aligned} \int_{\sigma(N)} |z - \lambda|^2 d\langle E(z)h, h \rangle &= \int_{\sigma(N)} |z - \lambda|^2 d\langle E(z)E(\{\lambda\})h, h \rangle \\ &= \int_{\{\lambda\}} |z - \lambda|^2 d\langle E(z)h, h \rangle = 0 \end{aligned}$$

That means

$$\|(N - \lambda)h\|^2 = \int_{\sigma(N)} |z - \lambda|^2 d\langle E(z)h, h \rangle = 0$$

Then $h \in \ker(N - \lambda)$, i.e. $\lambda \in \sigma_p(N)$.

2) and 3) can be obtained directly by above proof. \square

Remark. This means that if λ is an isolated point in $\sigma(N)$, then $\lambda \in \sigma_p(N)$.

Now, we can extend the **Example 4** in the subsection **2.2.2**.

Example 3. Let μ be a compactly supported, regular Borel measure on \mathbb{C} and (X, Ω, μ) be the measure space. For each $\phi \in L^\infty(\mu)$, we define the map

$$\begin{aligned} M_\phi: L^2(\mu) &\longrightarrow L^2(\mu) \\ f(z) &\longmapsto \phi(z)f(z) \end{aligned}$$

- 1) the spectral measure for M_ϕ is $M_{\chi_{\phi^{-1}(\Delta)}}$.
- 2) If $\psi \in B(\sigma(M_\phi))$, then $\psi(M_\phi) = M_{\psi \circ \phi}$.

4.2 Abelian Von Neumann Algebras

In above section, we built the representation ρ from $B(X)$ to $\mathcal{B}(\mathcal{H})$, and more over, ρ is *wk* – *WOT* continuous. Thus we want to prove $\rho(B(X))$ is indeed a *WOT*-closed C^* -subalgebra of $\mathcal{B}(\mathcal{H})$. If $X = \sigma(N)$, $\rho(B(\sigma(N)))$ is the *WOT*-closed C^* -subalgebra generated by N , denoted by $W^*(N)$. We will provide more rigorous proof of this statement. After that, we will make the structure of $W^*(N)$ more explicit. But before doing these, we some properties of *WOT*-closed subalgebra.

4.2.1 Double Commutant Theorem

Definition 1. A von Neumann algebra is a weakly (or strongly) closed C^* -subalgebra of $\mathcal{B}(\mathcal{H})$ that contains the identity.

Definition 2. If $\mathcal{S} \subset \mathcal{B}(\mathcal{H})$ is a subset, then the commutant of \mathcal{S} , denoted by \mathcal{S}' ,

$$\mathcal{S}' = \{ T \in \mathcal{B}(\mathcal{H}) : TS = ST \ \forall S \in \mathcal{S} \}$$

There are some easy properties if the commutant.

Proposition 3. Let \mathcal{S} be a subset of $\mathcal{B}(\mathcal{H})$.

- 1) \mathcal{S}' is a von Neumann algebra.
- 2) If $A = \oplus_{i=1}^n A_i$ acts on $\mathcal{H}^{(n)}$ and $T = [T_{ij}] \in \mathcal{B}(\mathcal{H}^{(n)})$, then $T \in \{A\}'$ if and only if $T_{ij}A_j = A_iT_{ij}$ for all $0 \leq i, j \leq n$.
- 3) $[\mathcal{S}^{(n)}]'' = [\mathcal{S}''^{(n)}]$.

Definition 4. Let \mathcal{S} be a linear manifold of $\mathcal{B}(\mathcal{H})$.

$$\text{Lat } \mathcal{S} = \{ \mathcal{M} : \mathcal{M} \text{ is a invariant closed subspace for } S, \ \forall S \in \mathcal{S} \}$$

Proposition 5. If \mathfrak{A} is a C^* -subalgebra of $\mathcal{B}(\mathcal{H})$ containing the identity, then

$$\overline{A}^{SOT} = \{ T \in \mathcal{B}(\mathcal{H}) : \text{Lat } \mathfrak{A}^n \subset \text{Lat } T^n \ \forall n \in \mathbb{N} \}$$

Proof. Assume $T \in \overline{A}^{SOT}$, there exist $T_k \in \mathfrak{A}$ s.t. $T_k h \rightarrow Th$ for all $h \in \mathcal{H}$. If $\mathcal{M} \in \text{Lat } \mathfrak{A}^{(n)}$, $T_k^{(n)}(\mathcal{M}) \subset \mathcal{M}$. Since \mathcal{M} is closed, $T^{(n)}(\mathcal{M}) \subset \mathcal{M}$. Conversely, let T satisfy the condition. Fix $h_1, \dots, h_n \in \mathcal{H}$ and $\varepsilon > 0$. It suffices to show that there exists $A \in \mathfrak{A}$ s.t. $\|(T - A)h_k\| < \varepsilon$ for $1 \leq k \leq n$. Let $\mathcal{M} = \overline{\text{span}\{(Ah_1, \dots, Ah_n) : A \in \mathfrak{A}\}}$. Since $1 \in \mathfrak{A}$, $(h_1, \dots, h_n) \in \mathcal{M}$. Clearly, $\mathcal{M} \in \text{Lat } \mathfrak{A}^{(n)}$. Therefore $\mathcal{M} \in \text{Lat } T^{(n)}$, i.e. $T^{(n)}(\mathcal{M}) \subset \mathcal{M}$

$$T^{(n)}(h_1, \dots, h_n) = (Th_1, \dots, Th_n) \in \mathcal{M}$$

Thus, clearly for any given $\varepsilon > 0$, there exists $A \in \mathfrak{A}$ s.t.

$$\|(Ah_1, \dots, Ah_n) - (Th_1, \dots, Th_n)\| < \varepsilon$$

That means,

$$\sum_{k=1}^n \|(A - T)h_k\|^2 < \varepsilon^2$$

Therefore, $\|(T - A)h_k\| < \varepsilon$ for $1 \leq k \leq n$. \square

Theorem 6 (Double Commutant Theorem). *If \mathfrak{A} is a C^* -subalgebra of $\mathcal{B}(\mathcal{H})$ containing the identity, then*

$$\overline{\mathfrak{A}}^{SOT} = \overline{\mathfrak{A}}^{WOT} = \mathfrak{A}''$$

Proof. Firstly, by above proposition, \mathfrak{A}'' is SOT -closed and clearly $\mathfrak{A} \subset \mathfrak{A}''$. Therefore, it is sufficient to show that $\mathfrak{A}'' \subset \overline{\mathfrak{A}}^{SOT}$. Let $T \in \mathfrak{A}''$. Then, by above proposition, we just to check that $\text{Lat } \mathfrak{A}^{(n)} \subset \text{Lat } T^{(n)}$. Let $\mathcal{M} \in \text{Lat } \mathfrak{A}^{(n)}$. Thus, \mathcal{M} reduces $\mathfrak{A}^{(n)}$. That means the projection P from $\mathcal{H}^{(n)}$ to \mathcal{M} satisfies $P\mathfrak{A}^{(n)} = \mathfrak{A}^{(n)}P$, i.e. $P \in \{\mathfrak{A}^{(n)}\}'$. Since $T \in \{\mathfrak{A}\}''$, and by above proposition $T^{(n)} \in \{\mathfrak{A}^{(n)}\}''$, $PT^{(n)} = T^{(n)}P$. That means $\mathcal{M} \in \text{Lat } T^{(n)}$ \square

Remark. Therefore, if \mathfrak{A} is a von Neumann algebra, $\mathfrak{A} = \mathfrak{A}''$. It supports the statement that the algebraic structure in a C^* -algebra can completely determine the topological structure.

4.2.2 Abelian Von Neumann Algebras

Proposition 1. 1) *If $\{\mathfrak{A}_i\}_{i \in I}$ is a collection of von Neumann algebras, then $\cap_{i \in I} \mathfrak{A}_i$ is a von Neumann algebra.*

2) *If \mathfrak{A} is von Neumann algebra, then the center of \mathfrak{A} , $Z_{\mathfrak{A}} = \mathfrak{A} \cap \mathfrak{A}'$ is a von Neumann algebra.*

We said that the norm closed C^* -subalgebras of $\mathcal{B}(\mathcal{H})$ not containing enough interesting elements. Now, we can see what are these "interesting elements" the WOT -closed C^* -subalgebras have.

Theorem 2. *Let \mathfrak{A} be a von Neumann algebra, then \mathfrak{A} is the norm closed linear span of the projections in \mathfrak{A} .*

Proof. Since \mathfrak{A} is the linear span of $\text{Re } \mathfrak{A}$, we just need to show $\text{Re } \mathfrak{A}$ is the norm closed linear span of the projections in \mathfrak{A} .

Let $A \in \text{Re } \mathfrak{A}$ with the spectral measure E . E is defined on $\sigma(A)$. If $[a, b] \subset \sigma(A)$, then there exists a sequence $\{u_n(t)\}$ in $C(\sigma(A))$ s.t.

$$u_n(t) = 1 \text{ for } a \leq t \leq b - n^{-1} \text{ and } u_n(t) = 0 \text{ for other } t$$

Then $u_n(t) \rightarrow \chi_{[a,b]}(t)$. If $h \in \mathcal{H}$,

$$\|(u_n(A) - E[a,b])h\|^2 = \int_{\sigma(A)} |u_n(t) - \chi_{[a,b]}(t)|^2 dE_{hh}(t)$$

Therefore, $u_n(A) \rightarrow E[a,b]$ in *SOT*. And since $u_n(A) \in \mathfrak{A}$ and \mathfrak{A} is *SOT*-closed, $E[a,b]$ is in \mathfrak{A} . Thus, for any Borel subset Ω of $\sigma(A)$, $E(\Omega) \in \mathfrak{A}$. Then $A = \int_{\sigma(A)} t dE$ is the norm closed of linear span of all $E(\Omega)$ \square

Remark. In fact, for any normal operator N in a von Neumann algebra \mathfrak{A} with the spectral measure E and for any $\phi \in B(\sigma(N))$, $\phi(N)$ is in \mathfrak{A} . It is because that $\overline{C(\sigma(N))}^{WOT} = B(\sigma(N))$, and the map from $B(\sigma(N))$ to $\mathcal{B}(\mathcal{H})$ is *wk-WOT* continuous. That means $\phi(N)$ can be approximated by a sequence in \mathfrak{A} with respect to the *WOT*, thus $\phi(N)$ is in \mathfrak{A} . In particular, for any Borel subset Ω of $\sigma(N)$, the projection $E(\Omega) \in \mathfrak{A}$.

So, these interesting elements are projections. We say the von Neumann algebra have enough projections which can generate the von Neumann algebra. Projections play an important role in a von Neumann algebra. We will classify von Neumann algebras by classifying the containing projections.

Proposition 3. *Let \mathfrak{A} be a von Neumann algebra and $A \in \mathfrak{A}$.*

- 1) *If E and F are the projections onto $\overline{\text{ran } A}$ and $\ker A$, then $E, F \in \mathfrak{A}$.*
- 2) *If $A = W|A|$ is the polar decomposition, then $|A|, W \in \mathfrak{A}$.*

Proof. For 1), if $h \in \ker A$ and $B \in \mathfrak{A}'$, then $BAh = ABh = 0$. That means F is invariant for $\{\mathfrak{A}'\}'$ and thus $F \in \{\mathfrak{A}\}''$. And $E = \ker A^* \in \{\mathfrak{A}\}''$.

For 2), $|A|$ is clearly in \mathfrak{A} by the Continuous Functional Calculus. For W , we just need to prove that $WT = TW$ for any $T \in \mathfrak{A}'$. Let $E = \ker |A|$. By 1), $E \in \mathfrak{A} = \mathfrak{A}''$. For $x \in \ker |A| = E$, $ET = TE$ implies $Tx \in \ker |A|$. Since $\ker |A| = \ker A$, $WTx = 0 = TWx$. On the other hand, if $x \in \text{ran } |A|$, $x = |A|y$ for some $y \in \mathcal{H}$.

$$WTx = WT|A|y = W|A|Ty = ATy = T Ay = TW|A|y = TWx$$

And this equation can be easily extended to $\overline{\text{ran } |A|}$. $\mathcal{H} = \ker |A| \oplus \overline{\text{ran } |A|}$, thus $TW = WT$. \square

Before researching the structure of abelian von Neumann algebras, there is a typical example.

Example 4. If (X, Ω, μ) is a finite measure space, and $\mathfrak{A}_\mu = \{M_\phi : \phi \in L^\infty(\mu)\}$, then

$$\mathfrak{A}_\mu = \mathfrak{A}'_\mu = \mathfrak{A}''_\mu$$

Proof. It suffices to show $\mathfrak{A}_\mu = \mathfrak{A}'_\mu$. And since \mathfrak{A}_μ is abelian, $\mathfrak{A}_\mu \subset \mathfrak{A}'_\mu$. Consider $T \in \mathfrak{A}'_\mu$, then $\mu(X) < \infty$ implies that $1 \in L^2(\mu)$. Let $\phi = T(1) \in L^2(\mu)$.

Check: $T = M_\phi$

Firstly, if $\psi \in L^\infty(\mu)$,

$$T(\psi) = T(M_\psi 1) = M_\psi T(1) = \phi \psi$$

And also,

$$\|T(\psi)\| = \|\phi \psi\| \leq \|T\| \|\psi\|_2$$

Using this we can show $\phi \in L^\infty(\mu)$. Let

$$\Delta_n = \{ x \in X : |\phi(x)| \leq n \}, \quad \psi = \chi_{\Delta_n}$$

Then we have

$$\|T\|^2 \mu(\Delta_n) = \|T\|^2 \|\psi\|^2 \geq \|\phi\|^2 \|\psi\|^2 = \int_{\Delta_n} |\phi|^2 d\mu = n^2 \mu(\Delta_n)$$

Thus $\mu(\Delta_n) = 0$. $\phi \in L^\infty(\mu)$. And since $L^\infty(\mu)$ is dense in $L^2(\mu)$, $T = M_\phi$ can be defined on $L^2(\mu)$. \square

Remark. In fact, this result can be also true for (X, Ω, μ) σ -finite.

Definition 5. Let $\mathcal{S} \subset \mathcal{B}(\mathcal{H})$ and $e_0 \in \mathcal{H}$.

- 1) e_0 is said to be star-cyclic for \mathcal{S} if $\overline{C^*(\mathcal{S})e_0} = \mathcal{H}$.
- 2) e_0 is said to be separating for \mathcal{S} there is no nonzero $S \in \mathcal{S}$ s.t. $Se_0 = 0$.

Proposition 6. If \mathfrak{A} is a unital C^* -algebra of operators, a vector is cyclic if and only if it is separating for \mathfrak{A}'

Proof. For $T \in \mathfrak{A}'$, s.t. $Te_0 = 0$. Then $T\mathfrak{A}e_0 = \mathfrak{A}Te_0 = 0$. If e_0 is the cyclic vector for \mathfrak{A} , then $\mathfrak{A}e_0$ is dense in \mathcal{H} . Therefore, $T = 0$.

Conversely, if e_0 separates \mathfrak{A}' , let P be the projection onto $\overline{(\mathfrak{A}e_0)}^\perp$. Then $P \in \mathfrak{A}'$ and $Pe_0 = 0$. Therefore, $P = 0$. e_0 is cyclic for \mathfrak{A} . \square

Corollary 7. If \mathfrak{A} is an abelian operator algebra, then each cyclic vector for \mathfrak{A} is also a separating vector for \mathfrak{A} .

Definition 8. An abelian subalgebra $\mathfrak{A} \subset \mathcal{B}(\mathcal{H})$ is said to be a maximal abelian subalgebra if there is no abelian subalgebra in $\mathcal{B}(\mathcal{H})$ containing \mathfrak{A} . In particular, a maximal abelian von Neumann algebra is called a masa.

Theorem 9. Let \mathfrak{A} be a subalgebra in $\mathcal{B}(\mathcal{H})$. Then the following statement are equivalent.

1) \mathfrak{A} is a masa.

2) $\mathfrak{A} = \mathfrak{A}'$

Proof. Assume \mathfrak{A} is a masa. Then $\mathfrak{A} \subset \mathfrak{A}'$. If there is a nonzero $A \in \mathfrak{A}' \setminus \mathfrak{A}$, then $\mathbb{C} = W^*(\mathfrak{A}, A)$ is an abelian von Neumann algebra containing \mathfrak{A} , which is a contradiction. Thus $\mathfrak{A} = \mathfrak{A}'$.

Conversely, \mathfrak{A} is an abelian von Neumann algebra by above proposition. By Zorn's Lemma there is a maximal abelian von Neumann algebra \mathfrak{B} s.t. $\mathfrak{A} \subset \mathfrak{B}$. Then $\mathfrak{B}' \subset \mathfrak{A}'$. By above statement, $\mathfrak{B} = \mathfrak{B}'$. Thus $\mathfrak{B} \subset \mathfrak{A}$. \square

Theorem 10. *If \mathcal{H} is separable and \mathfrak{A} is an abelian C^* -algebra containing the identity, the following statements are equivalent.*

1) $\mathfrak{A} = \mathfrak{A}'$.

2) \mathfrak{A} is strongly closed and has a cyclic vector.

3) There is a compact metric space X , a regular Borel measure μ with support X , and a unitary $U: L^2(\mu) \rightarrow \mathcal{H}$ s.t. $U\mathfrak{A}_\mu U^{-1} = \mathfrak{A}$.

Proof. 1) \Rightarrow 2) Let $\{e_n\}_{n=1}^\infty$ be a sequence of unit vectors s.t.

$$\overline{\mathfrak{A}e_n} \perp \overline{\mathfrak{A}e_m}, \text{ for } n \neq m \text{ and } \mathcal{H} = \bigoplus_{n=1}^\infty \overline{\mathfrak{A}e_n}$$

Set $e_0 = \sum_{n=1}^\infty 2^{-n} e_n$ and $P_n = \overline{\mathfrak{A}e_n}$. Then $P_n \in \mathfrak{A}' = \mathfrak{A}$.

$$e_n = 2^n P_n e_0 \in \overline{\mathfrak{A}e_0} \Rightarrow \overline{\mathfrak{A}e_n} \subset \overline{\mathfrak{A}e_0}$$

Therefore, $\mathcal{H} = \overline{\mathfrak{A}e_0}$.

2) \Rightarrow 3) Since \mathcal{H} is separable, the unit ball in \mathfrak{A} is a WOT -closed compact metric space, therefore there is a countable weakly dense set. Let \mathfrak{A}_1 be the C^* -algebra generated by this set and X be the maximal ideal space of \mathfrak{A}_1 . Clearly, X is metrizable. Then the inverse of Gelfand transform

$$\rho: C(X) \longrightarrow \mathfrak{A}_1 \subset \mathfrak{A}$$

And there is a spectral measure E , s.t. $\rho(f) = \int_X f dE$. Let e_0 be the cyclic vector of \mathfrak{A} . Then we can define the U ,

$$\begin{aligned} U: C(X) &\longrightarrow \mathfrak{A}_1 e_0 \\ \phi &\longmapsto \rho(\phi) e_0 \end{aligned}$$

Let $\mu(\Delta) = \langle E(\Delta) e_0, e_0 \rangle$, then $C(X) \subset L^2(\mu)$ by the fact that μ is positive finite measure. For $\phi \in C(X)$,

$$\|\rho(\phi) e_0\|^2 = \langle \rho(\phi^*) \rho(\phi) e_0, e_0 \rangle = \int |\phi|^2 d\mu = \|\phi\|^2$$

Therefore, U is an isometry. Since $C(X)$ is norm dense in $L^2(\mu)$ and $\overline{\mathfrak{A}_1 e_0} = \overline{\mathfrak{A} e_0} = \mathcal{H}$, $U: L^2(\mu) \rightarrow \mathcal{H}$ is a unitary. Let $\psi \in C(X)$ and $\phi \in C(X)$

$$UM_\psi \phi = U(\psi \phi) = \rho(\psi \phi) e_0 = \rho(\psi) \rho(\phi) e_0 = \rho(\psi) U \phi$$

Since $\rho(C(X)) = \mathfrak{A}_1$, $U^* \mathfrak{A}_1 U = \{M_\phi: \phi \in C(X)\}$. Define the map $\Gamma: \mathfrak{A} \rightarrow \mathfrak{A}_\mu$ by $\Gamma(A) = U^* A U$. We have known that

$$\Gamma(\mathfrak{A}_1) \subset \{M_\phi: \phi \in C(X)\} \subset \mathfrak{A}_\mu,$$

$$\overline{\{M_\phi: \phi \in C(X)\}}^{WOT} \subset \mathfrak{A}_\mu$$

And Γ is clearly WOT -closed, thus

$$\Gamma(\mathfrak{A}) \subset \overline{\Gamma(\mathfrak{A}_1)}^{WOT} \subset \mathfrak{A}_\mu$$

Γ is well-defined. In fact, by similar proof of $\mathfrak{A}_\mu = \mathfrak{A}'_\mu$, $\{M_\phi: \phi \in C(X)\}' = \mathfrak{A}_\mu$, i.e.

$$\mathfrak{A}_\mu = \overline{\{M_\phi: \phi \in C(X)\}}^{WOT}$$

Finally, it can easily check $\overline{\Gamma(\mathfrak{A}_1)}^{WOT} \subset \Gamma(\mathfrak{A})$. Therefore, Γ is an isomorphism.

3) \Rightarrow 1) It clearly holds by above example. \square

This theorem gives an explicit structure of the maximal abelian von Neumann algebra, which is equivalent to $L^\infty(\mu)$. And for a general abelian von Neumann algebra, by Zorn's Lemma, it can be embedded into a masa.

Corollary 11. *If \mathcal{H} is separable, every abelian C^* -algebra contained in $\mathcal{B}(\mathcal{H})$ has a separating vector.*

Proof. By Zorn's Lemma, there is a masa containing the abelian C^* -algebra and thus by above theorem it has a cyclic vector, which can separate the masa, in particular, this C^* -algebra. \square

4.2.3 Star-Cyclic Normal Operators

In above subsection, we can see $\mathfrak{A}_\mu = \{M_\phi: \phi \in L^\infty(\mu)\} \cong L^\infty(\mu)$ play an important role in the abelian von Neumann algebra. So we will see more properties of \mathfrak{A}_μ .

Theorem 1 (Fuglede-Putnam Theorem). *Let N and M be normal operators and $A \in \mathcal{B}(\mathcal{H})$. If $NA = AM$, then $N^*A = AM^*$.*

Proof. Since for any polynomial, $p(N)A = Ap(M)$, then for a fixed $\omega \in \mathbb{C}$, by taking limits

$$e^{i\omega N}A = Ae^{i\omega M}, \text{ or } A = e^{-i\omega N}Ae^{i\omega M}$$

By the fact that $e^Xe^Y = e^{X+Y}$ for commuting operators X, Y , therefore

$$\begin{aligned} f(\omega) &= e^{-i\omega N^*}Ae^{i\omega M^*} \\ &= e^{-i\omega N^*}e^{-i\bar{\omega}N}Ae^{i\bar{\omega}M}e^{i\omega M^*} \\ &= e^{-i(\omega N^* + \bar{\omega}N)}Ae^{i(\omega M^* + \bar{\omega}M)} \end{aligned}$$

Note that f is a entire operator-valued function. Since $(\omega N^* + \bar{\omega}N)$ and $(\omega M^* + \bar{\omega}M)$ are self-adjoint, $e^{-i(\omega N^* + \bar{\omega}N)}$ and $e^{i(\omega M^* + \bar{\omega}M)}$ are unitaries. Then f is bounded, thus f is constant, i.e.

$$0 = f'(\omega) = -iN^*e^{-i\omega N^*}Ae^{i\omega M^*} + ie^{-i\omega N^*}AM^*e^{i\omega M^*}$$

Let $\omega = 0$, then $0 = -iN^*A + iAM^*$. \square

Corollary 2. *If N is a normal operator with $N = \int zdE$, then $NA = AN$ if and only if $E(\Delta)N = NE(\Delta)$ for all Borel set Δ .*

Then using this theorem we can find the generator of \mathfrak{A}_μ .

Theorem 3. *If μ is any compactly supported measure on \mathbb{C} , then*

$$\{N_\mu\}' = \mathfrak{A}'_\mu$$

Proof. Clearly, $\mathfrak{A}'_\mu \subset N'_\mu$. If $T \in N'_\mu$, then $TN_\mu^* = N_\mu^*T$. Therefore, T can commutes with any polynomials in z and \bar{z} , and taking limits, $T \in \mathfrak{A}'_\mu$. \square

Remark. By this, we can also see the statement in above theorem

$$\{M_\phi: \phi \in C(X)\}' = \mathfrak{A}_\mu$$

is true. Moreover, $\{N_\mu\}' = \mathfrak{A}'_\mu$ implies that

$$W^*(N_\mu) = \mathfrak{A}_\mu \cong L^\infty(\mu)$$

Therefore, in order to research \mathfrak{A}_μ , we just research N_μ .

We have seen that if μ is finite then N_μ is indeed a cyclic normal operator on $L^2(\mu)$ as $\overline{C^*(\mu)1} = L^2(\mu)$ by the Stone-Weierstrass Theorem. Now, we can prove any cyclic normal operator is equivalent to a N_μ for some finite μ .

Theorem 4. *A normal operator N is cyclic if and only if it is equivalent to N_μ for some compactly supported measure μ on \mathbb{C} . Moreover, if e_0 is the cyclic vector, then the unitary $U: \mathcal{H} \rightarrow L^2(\mu)$ with $N_\mu = UNU^*$ can be chosen s.t. $Ue_0 = 1$.*

Proof. If E is the spectral measure s.t. $N = \int_{\sigma(N)} z dE$, for any Borel set Δ in $\sigma(N)$, then define

$$\mu(\Delta) = \langle E(\Delta)e_0, e_0 \rangle = \|E(\Delta)e_0\|^2$$

μ is finite, positive, and compactly supported on \mathbb{C} . Then $C(\sigma(N)) \subset L^2(\mu)$. We know that $C^*(N) \cong C(\sigma(N))$. Define

$$\begin{aligned} U: \quad C^*(N)e_0 &\longrightarrow C(\sigma(N)) \\ \phi(N)e_0 &\longmapsto \phi \end{aligned}$$

Then we have,

$$\|\phi(N)e_0\|^2 = \langle \phi(N^*)\phi(N)e_0, e_0 \rangle = \int_{\sigma(N)} |\phi|^2 d\mu = \|\phi\|_2^2$$

U is a unitary. By taking norm limits, $U: \mathcal{H} \rightarrow L^2(\mu)$ is a unitary and $Ue_0 = 1$ and $UNe_0 = z$. If $\phi \in C(\sigma(N))$, then

$$UNU^*(\phi) = UN\phi(N)e_0 = U(z\phi)(N)e_0 = z\phi = N_\mu(\phi)$$

Also, taking the norm limits, $UNU^* = N_\mu$. \square

Theorem 5. *Two normal operator N_μ and N_ν are equivalent if and only if μ and ν are mutually absolutely continuous.*

Proof. If μ and ν are mutually absolutely continuous and $\phi = \frac{d\mu}{d\nu}$, then define

$$\begin{aligned} U: \quad L^2(\mu) &\longrightarrow L^2(\nu) \\ f &\longmapsto f\sqrt{\phi} \end{aligned}$$

Then we have

$$\|Uf\|^2 = \int |f|^2 \phi d\nu = \int |f|^2 d\mu = \|f\|^2$$

Thus U is an isometry. Since $\frac{1}{\phi} = \frac{d\nu}{d\mu}$, the map $g \rightarrow \sqrt{\frac{1}{\phi}}g$ is the inverse of U . Therefore, U is unitary, and

$$UN_\mu U^*g = UN_\mu(\sqrt{\frac{1}{\phi}}g) = U(z\sqrt{\frac{1}{\phi}}g) = zg = N_\nu g$$

Conversely, $N_\mu \cong N_\nu$ means $K = \text{supp } \mu = \sigma(N_\mu) = \sigma(N_\nu) = \text{supp } N_\nu$. Let $V: L^2(\mu) \rightarrow L^2(\nu)$ be the unitary s.t. $VN_\mu V^* = N_\nu$. Then also $VN_\mu^* V^* = N_\nu^*$. Therefore, for any polynomial in z and \bar{z} , let $V(1) = \psi$

$$Vp(N_\mu, N_\mu^*)(1) = p(N_\nu, N_\nu^*)V(1)$$

That means

$$\int_K |p|^2 d\mu = \int_K |p|^2 |\psi|^2 d\nu$$

Taking limits, for any $f \in C(K)$, above identity is also true. Thus, $d\mu = |\psi|^2 d\nu$, i.e. $\mu \ll \nu$. Similarly, by V^{-1} , we can get the converse. \square

4.3 The structure of $W^*(N)$

For a norm operator N , by above subsection, clearly

$$\{ M_\phi : \phi \in B \} \subset W^*(N) = \{N\}'$$

And we have seen, $\{N_\mu\}' = \mathfrak{A}_\mu \cong L^\infty(\mu)$. Firstly, we want to prove $\{M_\phi : \phi \in B\} = W^*(N)$. More, if we also can find some μ s.t. $W^*(N) = W^*(N_\mu)$, then the structure of $W^*(N)$ is explicit,

$$W^*(N) = \{ M_\phi : \phi \in B \} \cong L^\infty(\mu)$$

Let N be a normal operator acting on \mathcal{H} with the spectral measure E . If $h \in \mathcal{H}$ is nonzero, then let

$$\mu_h(\Delta) = E_{hh}(\Delta), \quad \mathcal{H}_h = \overline{W^*(N)h}, \quad N_h = N|_{\mathcal{H}_h}$$

In fact, \mathcal{H}_h is the smallest reducing space for N containing h . And N_h is a star-cyclic normal operator with the cyclic vector h . If E_h is the spectral measure of N_h , then by the uniqueness of spectral measure

$$E_h(\Delta) = \chi_\Delta(N_h) = \chi_\Delta(N)|_{\mathcal{H}_h} = E(\Delta)|_{\mathcal{H}_h}$$

Therefore, there is a unique unitary $U_h : \mathcal{H}_h \rightarrow L^2(\mu_h)$ s.t. $U_h h = 1$ and $U_h N_h U_h^* = N_{\mu_h}$.

Proposition 1. *Let N be a normal operator acting on \mathcal{H} and a nonzero $h \in \mathcal{H}$.*

$$\begin{array}{ccc} \rho_h : & W^*(N) & \longrightarrow & W^*(N_h) \\ & A & \longmapsto & A|_{\mathcal{H}_h} \end{array}$$

- 1) ρ_h is a $*$ -homomorphism and WOT-continuous.
- 2) If $\phi \in B(\sigma(N))$, then $\rho_h(\phi(N)) = \phi(N_h)$.
- 3) If h is separating for $W^*(N)$, then ρ_h is injective.
- 4) If $A \in W^*(N)$, then there exists a $\phi \in B(\sigma(N))$ s.t. $\rho_h(A) = \phi(N_h)$, and thus ρ_h is surjective.

Proof. For 1), if $p(z, \bar{z})$ is any polynomial, then $\rho(p(N, N^*)) = p(N_h, N_h^*)$. For any $A \in W^*(N)$, there exists $p_n(N, N^*)$ s.t. $p_n(N, N^*) \rightarrow A$ in WOT. Then for any $g, h \in \mathcal{H}_h$

$$\langle p_n(N, N^*)g, h \rangle \rightarrow \langle Ag, h \rangle$$

Therefore, $p_n(N_h, N_h^*) \rightarrow A|_{\mathcal{H}_h}$ in WOT. And thus, ρ_h is WOT-continuous and $\rho_h(A) \in W^*(N_h)$, well-defined. And clearly, ρ_h is a $*$ -homomorphism.

For 2), if $\phi \in B(\sigma(N))$, then there exist polynomials $p_n(z, \bar{z})$ s.t. $p_n \rightarrow \phi$ weakly in $M(\sigma(N))$ and $M(\sigma(N_h))$ as $\sigma(N_h) \subset \sigma(N)$. That means $p_n(N, N^*) \rightarrow \phi(N)$ in WOT and $p_n(N_h, N_h^*) \rightarrow \phi(N_h)$. But $\rho_h(p_n(N, N^*) = p_n(N_h, N_h^*)$ and ρ is WOT -continuous, $\rho_h(\phi(N)) = \phi(N_h)$.

For 3), if h is separating for $W^*(N)$, $A_h = 0$ i.e $AW^*(N)h = 0$. That means $AW^*(N) = 0$, then $A = 0$.

For 4), if $A \in W^*(N)$ and $A_h = \rho_h(A)$, then $A_h \in \{N_h\}' = U_h^* \{N_{\mu_h}\}' U_h$. There is a $\psi \in L^\infty(\mu_h)$ s.t. $A_h = U_h^* M_\psi U_h$. Then there exist $\psi_n \in B(\sigma(N))$ s.t. $M_{\psi_n} \rightarrow M_\psi$ in norm. Thus $\psi_n(N_h) \rightarrow A_h$ in norm, there is a $\phi \in M(\sigma(N))$, s.t. $\phi(N_h) = A_h$. Moreover, this also implies ρ_h is surjective. \square

Remark. The most important part in this proof is that for the μ_h ,

$$W^*(N_{\mu_h}) = \mathfrak{A}_{\mu_h} = \{ M_\phi : \phi \in L^\infty(\mu_h) \} = \{ M_\phi : \phi \in B(\sigma(N_{\mu_h})) \}$$

Using above proposition we get one of the main result of this subsection.

Theorem 2. *If N is a normal operator, then*

$$W^*(N) = \{ \phi(N) : \phi \in B(\sigma(N)) \}$$

Proof. It is sufficient to show that $\mathfrak{A} = \{ \phi(N) : \phi \in B(\sigma(N)) \}$ is WOT -closed. Let $A \in \overline{\mathfrak{A}}^{WOT} = W^*(N)$. Then there exists $\{\phi_n\} \subset B(\sigma(N))$ s.t. $\phi_n(N) \rightarrow A$ in WOT . And for any $h \in \mathcal{H}$, there exists $\phi_h \in B(\sigma(N_h))$, s.t. $A_h = \phi_h(N_h)$, also $\phi_n(N_h) \rightarrow A_h$ in WOT . That means

$$Ah = \phi_h(N_h)h = \phi_h(N)h, \quad \forall h \in \mathcal{H}$$

But ϕ_h is dependent with the chose of h . So we want to find a $e \in \mathcal{H}$ s.t. $\phi_h(N_h) = \phi_e(N_h)$ for any $h \in \mathcal{H}$. Fix a $e \in \mathcal{H}$

Check: $\phi_h(N_h)h = \phi_e(N_h)h$

By the definition, for any Borel subset Δ of $\sigma(N)$, since

$$\begin{aligned} \mu_e &= \langle E(\Delta)e, e \rangle = \|E(\Delta)e\|^2 \\ \mu_h &= \langle E(\Delta)h, h \rangle = \|E(\Delta)h\|^2 \end{aligned}$$

μ_e and μ_h are mutually absolutely continuous. In fact, for any $g \in \mathcal{H}$, μ_h and μ_g are mutually absolutely continuous. Since $\phi_n(N_h) \rightarrow A_h$ in WOT , $\phi_n \rightarrow \phi_h$ weak* in $L^\infty(\mu_h)$. Also, $\phi_n \rightarrow \phi_e$ weak* in $L^\infty(\mu_e)$. Then for any Borel subset Δ

$$\int_{\Delta} \phi_n d\mu_h = \int_{\Delta} \phi_n \frac{d\mu_h}{d\mu_e} d\mu_e \rightarrow \int_{\Delta} \phi_e d\mu_h$$

and also $\int_{\Delta} \phi_n d\mu_h \rightarrow \int_{\Delta} \phi d\mu_h$, therefore, $\phi_h = \phi_e$ a.e. with respect to μ_h , in fact, it can be with respect to any μ_g . Then for $g \in \mathcal{H}$,

$$\langle \phi_h(N_h)g, g \rangle = \langle \phi_h(N)g, g \rangle = \int \phi_h d\mu_g = \int \phi_e d\mu_g = \langle \phi_e(N_h)g, g \rangle$$

Therefore, $\phi_h(N_h)h = \phi_e(N_h)h$. Then $A = \phi_e(N)$. \square

And also, by above proposition, the second result can also be proved.

Theorem 3. *If N is a normal operator on a separable Hilbert space, then there is a compactly supported, positive and finite measure μ , s.t. there is a map $\rho: L^\infty \rightarrow W^*(N)$, s.t.*

1) ρ is an isometrical $*$ -isomorphism.

2) $\rho: (L^\infty(\mu), \text{weak}^*) \rightarrow (W^*(N), WOT)$ is a homeomorphism.

Proof. Since $W^*(N)$ is an abelian von Neumann algebra, it has a separating vector e by the **Corollary 11** in the subsection **3.2.2**. Then let $\rho' = \rho_e$ and $\mu = \mu_e$. By above proposition, ρ' is a $*$ -isomorphism from $W^*(N_e)$ to $W^*(N)$. And since $W^*(N_e)$ is also $*$ -isomorphic to $L^\infty(\mu)$. Thus there is a $*$ -isomorphism ρ from $L^\infty(\mu)$ to $W^*(N)$ s.t. for any $\phi \in B(\sigma(N))$, $\rho(\phi) = \phi(N)$. And thus ρ is isometrical.

For 2), since $W^*(N) = \{\phi(N): \phi \in B(\sigma(N))\}$, if $\phi_n(N) \rightarrow \phi(N)$ in WOT , then we can easily know $\phi_n \rightarrow \phi$ in weak^* . And the converse is trivial. ρ is weak^* - WOT homeomorphism. \square

Now, we can see the structure of any abelian von Neumann algebra on a separable Hilbert space.

Theorem 4. *If \mathfrak{A} is an abelian von Neumann algebra on a separable Hilbert space \mathcal{H} , then there is a compactly supported, positive and finite measure μ s.t. $\mathfrak{A} \cong L^\infty(\mu)$, $*$ -isomorphism.*

Proof. Since \mathcal{H} is separable, the closed unit ball in \mathfrak{A} is WOT -compactly metrizable, thus there is a countable WOT -dense set of \mathfrak{A} . Since \mathfrak{A} can be generated by projections in \mathfrak{A} , thus this WOT -dense set can be chosen as a set of projections $\{P_n\}$. Let $\mathfrak{C} = C^*(\{P_n\})$, then \mathfrak{C} is abelian. Thus there is a compact space X , which is totally disconnected, s.t. $\mathfrak{C} \cong C(X)$. And let $\rho: C(X) \rightarrow \mathfrak{C}$ be the inverse of Gelfand Transform. Then $p_n = \rho^{-1}(P_n)$ is the characteristic function and continuous. Define

$$f = \sum_{n=1}^{\infty} 3^{-n}(2p_n - 1)$$

Claim: $C^*(f) = C(X)$

By the Stone-Weierstrass Theorem, we just need to show that f can separate the points of X . Since $\{p_n\}$ generates $C(X)$, $\{p_n\}$ can separate the points in X . Thus, for $x \neq y \in X$, there is a p_k , s.t. $p_k(x) \neq p_k(y)$ and choosing the

k be the smallest one satisfying this condition, then

$$\begin{aligned} |f(x) - f(y)| &= 2 \left| \sum_{n=k}^{\infty} 3^{-n} (p_n(x) - p_n(y)) \right| \\ &\geq 23^{-k} - 2 \sum_{n=k+1}^{\infty} 3^{-n} \\ &= 3^{-k} \end{aligned}$$

Therefore, let $A = \rho(f) = \sum_{n=1}^{\infty} 3^{-n} (2P_n - 1)$, $C^*(A) = \mathfrak{C}$. Moreover, $W^*(A) = \overline{C^*(A)}^{WOT} = \mathfrak{A}$. And since A is self-adjoint, thus by above theorem, there exists μ on $\sigma(A)$, s.t.

$$\mathfrak{A} \cong L^\infty(\mu) \quad \square$$

4.4 Multiplicity Theory

Finally, we want find the necessary and sufficient condition making two normal operators be equivalent. We should assume the Hilbert space is always separable.

Proposition 1. *If N is a normal operator on \mathcal{H} and $e \in \mathcal{H}$, then there is a separating vector e_0 for $W^*(N)$ s.t. $e \in \overline{W^*(N)}$.*

Proof. Choosing any separating operator f_0 for $\overline{W^*(N)}$, let $\mathcal{M} = \overline{W^*(N)}f_0$, then $e = m_0 + g_0$ for $m_0 \in \mathcal{M}$ and $g_0 \in \mathcal{M}^\perp$. Now, let $\mathcal{G} = \overline{W^*(N)}g_0$ and $e_0 = f_0 + g_0$. Then $e \in \overline{W^*(N)}e_0$.

Check: e_0 is separating for $W^*(N)$

For any $A \in W^*(N)$, since both \mathcal{M} and \mathcal{G} reduce $W^*(N)$, if

$$Ae_0 = Af_0 + Ag_0 = 0$$

then $Af_0 = 0$. Because f_0 is separating for $W^*(N)$, $A = 0$. \square

Proposition 2. *Let μ be a compactly supported measure on \mathbb{C} and Δ be a Borel subset of $\text{supp } \mu$ and $\nu = \mu|_\Delta$. If $N = N_\mu \oplus N_\nu$ on $L^2(\mu) \oplus L^2(\nu)$ and $g \oplus h \in L^2(\mu) \oplus L^2(\nu)$ with $h \neq 0$ a.e. ν , then there is a $f \in L^2(\mu)$ s.t. $f \oplus h$ is a separating vector for $W^*(N)$ and $g \oplus h \in \overline{W^*(N)}(f \oplus h)$.*

Proof. Let $f(z) = g(z)$ for $z \in \Delta$ and $f(z) = 1$ for $x \in \Delta'$ where $\Delta' = \text{supp } \mu \setminus \Delta$. In fact,

$$\mathcal{H} = \overline{W^*(N)}(f \oplus h) = \{ \phi f \oplus \phi h : \phi \in L^\infty(\mu) \}$$

Since $\phi \chi_{\Delta'} \oplus 0 = \phi \chi_{\Delta'}(f \oplus h)$,

$$g \oplus h = f \oplus h - (1 - g)\chi_{\Delta'} \oplus 0 \in \mathcal{H}$$

If $\phi \in L^\infty(\mu)$ with $\phi f \oplus \phi h = 0$, then $\phi f = \phi h = 0$ a.e. μ and ν respectively. By definition, $\phi = 0$ a.e. μ . \square

Then we can decompose normal operators as a "decreasing" sequence of cyclic operators.

Theorem 3. *If N is a normal operator, then there is a sequence of measures $\{\mu_n\}$ on \mathbb{C} s.t. $\mu_{n+1} \ll \mu_n$ for all n and*

$$N \cong N_{\mu_1} \oplus N_{\mu_2} \oplus N_{\mu_3} \oplus \dots$$

Proof. Let e_1 be a separating vector for $W^*(N)$ and $\{f_n\}$ be the orthonormal basis of \mathcal{H} s.t. $f_1 = e_1$. Put

$$\mathcal{H}_1 = \overline{W^*(N)e_1}, \quad \mu_1(\Delta) = \|E(\Delta)\|^2, \quad N_2 = N|_{\mathcal{H}_1}$$

Then for $f'_2 =$ the projection of f_2 to \mathcal{H}_1^\perp , there is a separating vector e_2 for $W^*(N_2)$ s.t.

$$f'_2 \in \mathcal{H}_2 = \overline{W^*(N_2)e_2}$$

And clearly, if $\mu_2\Delta = \|E(\Delta)|_{\mathcal{H}_2}e_2\|$, $\mu_2 \ll \mu_1$. By induction and the Zorn's Lemma,

$$N \cong N_{\mu_1} \oplus N_{\mu_2} \oplus N_{\mu_3} \oplus \dots \quad \square$$

Therefore, if there is another normal operator $M \cong N_{\nu_1} \oplus N_{\nu_2} \oplus \dots$, we want to find the relation between the equivalence of N and M and the equivalence of the coincided μ_n and ν_n . But before that we need some lemmas.

Proposition 4. *If N_1 and N_2 are normal operators on \mathcal{H}_1 and \mathcal{H}_2 respectively and $X: \mathcal{H}_1 \rightarrow \mathcal{H}_2$ s.t. $XN_1 = N_2X$, then*

- 1) $\overline{\text{ran } X}$ reduces N_2 .
- 2) $\ker X$ reduces N_1 .
- 3) If $M_1 = N_1|_{(\ker X)^\perp}$ and $M_2 = N_2|_{\overline{\text{ran } X}}$, then $M_1 \cong M_2$.

Proof. For 1), if $f_1 \in \mathcal{H}_1$, $XN_1f_1 = N_2Xf_1 \in \text{ran } X$. Therefore, $\overline{\text{ran } X}$ is invariant N_2 . And also, by the Fuglede-Putnam Theorem, $XN_1^* = N_2^*X$, thus $\overline{\text{ran } X}$ is invariant N_2^* .

2) is similar as 1). For 3), it is sufficient to show this when $\ker X = \{0\}$ and $\text{ran } X$ is dense. By the Polar Decomposition, $X = UA$ for some unitary U . Since $X^*N_2^* = N_1^*X^*$, $X^*N_2 = N_1X^*$. Therefore $A^2 \in \{N_1\}'$, so $A \in \{N_1\}'$ by the functional calculus.

$$N_2X = N_2UA = XN_1 = UAN_1 = UN_1A$$

That means $N_2U = UN_1$ on $\overline{\text{ran } A} = \overline{\text{ran } X} = \mathcal{H}_2$, therefore $N_1 \cong N_2$. \square

Proposition 5. If N, A and B are normal operators on \mathcal{H} , \mathcal{H}_A and \mathcal{H}_B respectively and N is $*$ -cyclic and

$$N \oplus A \cong N \oplus B$$

then $A \cong B$.

Proof. Let $U: \mathcal{H} \oplus \mathcal{H}_A \rightarrow \mathcal{H} \oplus \mathcal{H}_B$, and U can be

$$U = \begin{pmatrix} U_{11} & U_{12} \\ U_{21} & U_{22} \end{pmatrix}$$

Then we have

$$\begin{pmatrix} U_{11}N & U_{12}A \\ U_{21}N & U_{22}A \end{pmatrix} = \begin{pmatrix} NU_{11} & NU_{12} \\ BU_{21} & BU_{22} \end{pmatrix}$$

$$\begin{pmatrix} U_{11}N^* & U_{12}A^* \\ U_{21}N^* & U_{22}A^* \end{pmatrix} = \begin{pmatrix} N^*U_{11} & N^*U_{12} \\ B^*U_{21} & B^*U_{22} \end{pmatrix}$$

Then by above proposition,

$$A|_{(\ker U_{22})^\perp} = B|_{(\ker U_{22}^*)^\perp}$$

For $\ker U_{22}$ and $\ker U_{22}^*$, put $\mathcal{M}_1 = U_{12}(\ker U_{22})$, $\mathcal{M}_2 = U_{21}^*(\ker U_{22}^*)$, also by the identity, $U^*U = UU^* = 1$,

$$A|_{\ker U_{22}} \cong N|_{\mathcal{M}_1}, \quad B|_{\ker U_{22}^*} \cong N|_{\mathcal{M}_2}$$

And then it easily check that $\mathcal{M}_1 = \mathcal{M}_2$. □

Theorem 6. For two normal operators N and M and the two sequences of measures $\{\mu_n\}$ and $\{\nu_n\}$ on \mathbb{C} s.t. $\mu_{n+1} \ll \mu_n$ and $\nu_{n+1} \ll \nu_n$ for all n s.t.

$$\begin{aligned} N &\cong N_{\mu_1} \oplus N_{\mu_2} \oplus \cdots \\ M &\cong N_{\nu_1} \oplus N_{\nu_2} \oplus \cdots \end{aligned}$$

then $N \cong M$ if and only if μ_n and ν_n are mutually absolutely continuous for all n .

Proof. Let $N \cong M$ with the unitary U . If e_1 is the separating vector for $W^*(N)$, then $f_1 = Ue_1$ is the separating vector for $W^*(M)$. If the μ_1 and ν_1 are the coincided measures, then μ_1 and ν_1 are mutually absolutely continuous, thus $N_{\mu_1} \cong N_{\nu_1}$. Let $\mathcal{H}_1 = \overline{W^*(N)e_1}$ and $\mathcal{H}'_1 = \overline{W^*(M)f_1}$. Then by above theorem,

$$N = N_{\mu_1} \oplus N|_{\mathcal{H}_1^\perp} \cong M = N_{\nu_1} \oplus M|_{\mathcal{H}'_1^\perp}$$

we know that $N|_{\mathcal{H}_1^\perp} \cong M|_{\mathcal{H}'_1^\perp}$. Therefore, by induction, we can get the result. □

Then we can get another formula of above theorem to help us build the Multiplicity Theorem.

Corollary 7. *If N is a normal operator with the coincided measure μ , then there is a decreasing sequence $\{\Delta_n\}$ of Borel subsets of $\sigma(N)$ s.t. $\Delta_1 = \sigma(N)$ and*

$$N \cong N_\mu \oplus N_{\mu|_{\Delta_2}} \oplus N_{\mu|_{\Delta_3}} \oplus \cdots$$

Also, if M is another normal operator with similar formula

$$M \cong N_\nu \oplus N_{\nu|_{\Sigma_2}} \oplus N_{\nu|_{\Sigma_3}} \oplus \cdots$$

then $N \cong M$ if and only if μ and ν are mutually absolutely continuous and $\mu(\Delta_n \setminus \Sigma_n) = 0 = \mu(\Sigma_n \setminus \Delta_n)$

Theorem 8. *If N is a normal operator, then there are mutually singular measures $\mu_\infty, \mu_1, \mu_2, \dots$ s.t.*

$$N \cong N_\infty^{(\infty)} \oplus N_{\mu_1} \oplus N_{\mu_2}^{(2)} \oplus N_{\mu_3}^{(3)} \oplus \cdots$$

If M is another normal operator with the corresponding measures ν_∞, ν_1, \dots , then $N \cong M$ if and only if μ_n and ν_n are mutually absolutely continuous for $1 \leq n \leq \infty$.

Proof. Let the μ and $\{\Delta_n\}$ be the sequence in above corollary. Put

$$\Sigma_\infty = \bigcap_{n=1}^{\infty} \Delta_n, \Sigma_n = \Delta_n \setminus \Delta_{n+1} \text{ and } \mu_n = \mu|_{\Sigma_n}$$

for $1 \leq n \leq \infty$ and $\nu_n = \mu|_{\Delta_n}$ for $1 \leq n < \infty$. Now,

$$\Delta_n = \Sigma_\infty \bigcup (\Delta_n \setminus \Delta_{n+1}) \bigcup (\Delta_{n+1} \setminus \Delta_{n+2}) \bigcup (\Delta_{n+2} \setminus \Delta_{n+3}) \bigcup \cdots$$

Therefore, $\nu_n = \mu_\infty + \mu_n + \mu_{n+1} + \mu_{n+2} + \cdots$ and $\{\mu_n\}$ are pairwise singular. Thus,

$$N_{\nu_n} = N_{\mu_\infty} \oplus N_{\mu_n} \oplus N_{\mu_{n+1}} \oplus \cdots$$

Then we have

$$\begin{aligned} N &= N_{\nu_1} \oplus N_{\nu_2} \oplus \cdots \\ &= N_\infty^{(\infty)} \oplus N_{\mu_1} \oplus N_{\mu_2}^{(2)} \oplus N_{\mu_3}^{(3)} \oplus \cdots \end{aligned}$$

□

Now, we can define the multiplicity function m_N for a normal operator N with the corresponding measure μ . As above preceding theorem,

$$N \cong N_\infty^{(\infty)} \oplus N_{\mu_1} \oplus N_{\mu_2}^{(2)} \oplus N_{\mu_3}^{(3)} \oplus \cdots$$

If there are disjoint Borel sets $\Delta_\infty, \Delta_1, \Delta_2, \dots$ of \mathbb{C} s.t. μ_n and $m\mu|_{\Delta_n}$ are mutually absolutely continuous for all n , then define $m_N: \mathbb{C} \rightarrow \mathbb{N} \cup \{\infty\}$ as

$$m_N = \infty\chi_{\Delta_\infty} + \chi_{\Delta_1} + 2\chi_{\Delta_2} + 3\chi_{\Delta_3} + \dots$$

Clearly, m_N is dependent with the choice of $\{\Delta_n\}$, but by the definition, m_N can be identified as an equivalent class with respect to μ . And moreover,

$$\Delta_n = \{ z: m_N(z) = n \}$$

Therefore, we can recover the $\{\Delta_n\}$ by the m_N . By combining with above theorem, we can classify the normal operators.

Theorem 9. *Two normal operators are equivalent if and only if they have same multiplicity function.*

Chapter 5

Compact Operators, Fredholm Theory and Perturbations

The compact operator is as similar as the finite dimensional operator with respect to the spectrums, the representations and so on. In fact, all compact operators can form the unique closed ideal in $\mathcal{B}(\mathcal{H})$. Moreover, it contains some important classes of operators, like the trace class, which provides an extra topology on $\mathcal{B}(\mathcal{H})$ and this topology plays an important role in the von Neumann algebra. In fact, the topology is coincided with the *WOT* in some cases, but it can give a new method to obtain some properties of general von Neumann algebras. For non-compact operators, there is a "weak version" of the compact operator, the Fredholm operator. By researching the Fredholm operators, it provides us an approach to find some "invariant variable" under the compact perturbations. Finally, we have known that finite dimensional normal operator can always be diagonalizable, but for infinite dimensional case, this statement is not true. But by some small compact perturbation, the normal operator can be diagonalizable, and this result has some applications of representation theory.

5.1 Spectrums

5.1.1 Elementary Properties

Definition 1. Let \mathcal{H} be a Hilbert space and $B_{\mathcal{H}}$ be the closed unit ball in \mathcal{H} and $T \in \mathcal{B}(\mathcal{H})$. T is compact if and only if the closure of $T(B_{\mathcal{H}})$ is compact. Then let $\mathcal{B}_0(\mathcal{H})$ denote the set of all compact operators.

Remark. Clearly, if T is finite rank, $T \in \mathcal{B}_0(\mathcal{H})$. Denote $\mathcal{B}_{00}(\mathcal{H})$ as the

set of all finite rank operators. Therefore, $\mathcal{B}_{00}(\mathcal{H}) \subset \mathcal{B}_0(\mathcal{H})$.

$$\mathcal{B}_{00}(\mathcal{H}) = \{ T \in \mathcal{B}(\mathcal{H}) : \text{ran } T \text{ is finite dimension} \}$$

Since \mathcal{H} is a complete metric space, the compactness of $\overline{T(B_{\mathcal{H}})}$ is equivalent to that any sequence in $\overline{T(B_{\mathcal{H}})}$ has a convergent subsequence with the convergent point in $\overline{T(B_{\mathcal{H}})}$. By this fact, we can get the following theorem.

Theorem 2. *Let \mathcal{H} be a Hilbert space.*

- 1) $\mathcal{B}_0(\mathcal{H})$ is a linear space.
- 2) $\mathcal{B}_0(\mathcal{H})$ is closed in norm.
- 3) $\mathcal{B}_0(\mathcal{H})$ is an ideal in $\mathcal{B}(\mathcal{H})$.
- 4) If $T \in \mathcal{B}_0(\mathcal{H})$, then $T^* \in \mathcal{B}_0(\mathcal{H})$.

Remark. Therefore, $\mathcal{B}_0(\mathcal{H})$ is a C^* -subalgebra, and moreover, $\mathcal{B}_0(\mathcal{H})$ is a closed ideal of $\mathcal{B}(\mathcal{H})$.

Definition 3. Let \mathcal{H} be a Hilbert space and $T \in \mathcal{B}(\mathcal{H})$. T is called completely continuous if for any sequence $\{x_n\} \in \mathcal{H}$ with $x_n \rightarrow x$ weakly, then $Tx_n \rightarrow Tx$ in norm.

Proposition 4. $T \in \mathcal{B}(\mathcal{H})$ is compact if and only if T is completely continuous.

Proof. Let T be a compact operator and $\{x_n\} \in \mathcal{H}$ with $x_n \rightarrow 0$ weakly. By the Principle of Uniform Boundedness, $M = \sup_n \|x_n\| < \infty$. Assuming $M \leq 1$,

$$\{Tx_n\} \subset \overline{T(B_{\mathcal{H}})}$$

By the fact that $\overline{T(B_{\mathcal{H}})}$ is compact, there is a subsequence $\{x_{n_k}\}$ s.t. $Tx_{n_k} \rightarrow y$ in norm. Since T is also weakly continuous, $Tx_n \rightarrow 0$ weakly. Thus $y = 0$. Then T is completely continuous.

Conversely, assume T is completely continuous.

Firstly, if \mathcal{H} is separable and since \mathcal{H} is reflexive, $B_{\mathcal{H}}$ is a weak-compact metric space. Therefore, for any $\{x_n\} \subset B_{\mathcal{H}}$, there is a subsequence $\{x_{n_k}\}$ and x s.t. $x_{n_k} \rightarrow x$ weakly. Therefore, $Tx_{n_k} \rightarrow x$ in norm, i.e. T is compact. For general case, let $\{x_n\} \subset B_{\mathcal{H}}$, then $\mathcal{H}_1 = \overline{\text{span}\{x_n\}}$ is separable. Therefore, since

$$T|_{\mathcal{H}_1}(\{x_n\}) = T(\{x_n\}) \subset T|_{\mathcal{H}_1}(B_{\mathcal{H}_1})$$

T is compact. □

By above theorem, we can see the power of compact operators. Intuitively, compact operators can "strengthen" the topology. They map the weakly convergent sequences be norm convergent sequence. How can they

do that? We have known the weak topology is agree with the norm topology on the finite dimensional space. So, compact operators may be the "extension" of finite rank operators. In fact, we can describe this property more rigorously.

Theorem 5. $T \in \mathcal{B}(\mathcal{H})$ is compact if and only if there is a sequence $\{T_n\} \subset \mathcal{B}_{00}(\mathcal{H})$ s.t. $T_n \rightarrow T$ in norm.

Proof. Assume T is compact. Therefore, $\overline{T(B_{\mathcal{H}})}$ is separable. Let $\{e_n\}_{n=1}^{\infty}$ be the basis of the dense subspace of $\overline{T(B_{\mathcal{H}})}$. Then define projections

$$P_n: \mathcal{H} \longrightarrow \text{span}\{e_1, e_2, \dots, e_n\}$$

and let $T_n = P_n T$. Clearly, $T_n \in \mathcal{B}_{00}(\mathcal{H})$. Now, clearly, $\|Th - T_n h\| \rightarrow 0$ for any $h \in \mathcal{H}$. Since $\overline{T(B_{\mathcal{H}})}$ is compact, for any given $\varepsilon > 0$, there are $h_1, \dots, h_m \in \mathcal{H}$ s.t.

$$T(B_{\mathcal{H}}) \subset \bigcup_{i=1}^m B_{\varepsilon}(Th_i)$$

where $B_{\varepsilon}(h_i)$ is the open ball centred at Th_i with radius ε . Therefore, for any $h \in B_{\mathcal{H}}$, choose h_j s.t. $\|Th - Th_j\| < \varepsilon$.

$$\begin{aligned} \|Th - T_n h\| &\leq \|Th - Th_j\| + \|Th_j - T_n h_j\| + \|P_n(Th_j - Th)\| \\ &\leq 2\|Th - Th_j\| + \|Th_j - T_n h_j\| \\ &< 2\varepsilon + \|Th_j - T_n h_j\| \end{aligned}$$

Since $\|Th_j - T_n h_j\| < \varepsilon$ for any h_j and $n > n_0$ for some n_0 ,

$$\|Th - T_n h\| < 3\varepsilon \text{ for } n > n_0$$

Then $T_n \rightarrow T$ uniformly, i.e. $T_n \rightarrow T$ in norm.

The converse is clearly since $\mathcal{B}_{00}(\mathcal{H}) \subset \mathcal{B}_0(\mathcal{H})$ and $\mathcal{B}_0(\mathcal{H})$ is norm closed. \square

Corollary 6. All projections $\{P_i\}_{i \in I}$ in $\mathcal{B}(\mathcal{H})$ form an approximate identity for the ideal $\mathcal{B}_0(\mathcal{H})$.

5.1.2 Spectrums of Compact Operators

The spectrum of a compact operator also has similar properties as a finite rank operator.

Theorem 1. If $T \in \mathcal{B}_0(\mathcal{H})$ and $\lambda \neq 0$ satisfying

$$\inf \{ \|(T - \lambda)h\| : \|h\| = 1 \} = 0$$

then $\lambda \in \sigma_p(T)$.

Proof. There is a sequence $\{h_n\}$ with $\|h_n\| = 1$ s.t. $\|(T - \lambda)h_n\| \rightarrow 0$. Since T is compact, there is a subsequence $\{h_{n_k}\}$ and a h_0 s.t. $\|Th_{n_k} - h_0\| \rightarrow 0$, therefore

$$h_{n_k} = \frac{1}{\lambda}((\lambda - T)h_{n_k} + Th_{n_k}) \rightarrow \frac{1}{\lambda}h_0$$

Then $\|\frac{1}{\lambda}h_0\| = 1 = |\frac{1}{\lambda}| \|h_0\|$, thus $\|h_0\| \neq 0$. Since $Th_{n_k} \rightarrow \frac{1}{\lambda}Th_0$ and $Th_{n_k} \rightarrow h_0$,

$$h_0 = \frac{1}{\lambda}Th_0, \text{ i.e. } Th_0 = \lambda h_0$$

□

In fact, we have known that

$$\sigma_{ap}(T) = \sigma_l(T) = \{ \lambda \in \mathbb{C} : \inf \{ \|(T - \lambda)h\| : \|h\| = 1 \} = 0 \}$$

and $\lambda \notin \sigma_{ap}(T)$ is equivalent to that $\ker(T - \lambda) = \{0\}$ and $\text{ran}(T - \lambda)$ is closed. Combining with the Riesz Theorem, any closed and bounded subset in a normed space is compact if and only if the normed space is finite dimensional, we have the following corollary.

Corollary 2. *If $T \in \mathcal{B}_0(\mathcal{H})$ and $\lambda \neq 0$ and $\ker(T - \lambda) = \{0\}$, then $\text{ran}(T - \lambda)$ is closed.*

Here is another important corollary, which says the point spectrum play an important role.

Corollary 3. *If $T \in \mathcal{B}_0(\mathcal{H})$, $\lambda \notin \sigma_p(T)$ with $\lambda \neq 0$ and $\bar{\lambda} \notin \sigma_p(T^*)$, then $\lambda \notin \sigma(T)$.*

Proof. By above theorem, $\lambda \notin \sigma_p(T)$ with $\lambda \neq 0$ means $\lambda \notin \sigma_l(T)$, thus $\ker(T - \lambda) = \{0\}$ and $\text{ran}(T - \lambda)$ is closed.

Similarly, for T^* , $\ker(T^* - \bar{\lambda}) = \{0\}$ and $\text{ran}(T^* - \bar{\lambda})$ is closed. Therefore,

$$\text{ran}(T - \lambda) = \overline{\text{ran}(T - \lambda)} = (\ker(T^* - \bar{\lambda}))^\perp = \mathcal{H}$$

Thus $T - \lambda$ is a bijection, and by the Inverse Mapping Theorem, $(T - \lambda)^{-1} \in \mathcal{B}(\mathcal{H})$, i.e. $\lambda \notin \sigma(T)$. □

Remark. In other words, for a compact operator T , if $\lambda \in \sigma(T)$ with $\lambda \neq 0$, then $\lambda \in \sigma_p(T)$ or $\lambda \in \sigma_p(T^*)$.

In fact, any nonzero point in $\sigma(T)$ for a compact operator T is isolated and moreover, it is an eigenvalue. To prove this, we need some lemmas.

Lemma 4. *If \mathcal{M} and \mathcal{N} are two closed linear subspaces of \mathcal{H} with $\mathcal{M} \subset \mathcal{N}$, then for any $\varepsilon > 0$, there exists a $y \in \mathcal{N}$ with $\|y\| = 1$ s.t. $\text{dist}(y, \mathcal{M}) \leq 1 - \varepsilon$.*

Proof. For $y \in \mathcal{N}$, define $\delta(y) = \text{dist}(y, \mathcal{M})$. Choosing $y_1 \in \mathcal{N} \setminus \mathcal{M}$, there exists a $x_0 \in \mathcal{M}$ s.t.

$$\delta(y_1) \leq \|x_0 - y_1\| \leq (1 + \varepsilon)\delta(y_1)$$

Let $y_2 = y_1 - x_0$, then

$$(1 + \varepsilon)\delta(y_2) = (1 + \varepsilon) \inf \{\|y_1 - x_0 - x\| : x \in \mathcal{M}\} = (1 + \varepsilon)\delta(y_1)$$

Thus $(1 + \varepsilon)\delta(y_2) > \|x_0 - y_1\| = \|y_2\|$. Let $y = \|y_2\|^{-1} y_2$, then $y_2 \in \mathcal{N}$.

$$\begin{aligned} \|y - x\| &= \left\| \|y_2\|^{-1} y_2 - x \right\| = \|y_2\|^{-1} \|y_2 - \|y_2\| x\| \\ &> ((1 + \varepsilon)\delta(y_2))^{-1} \|y_2 - \|y_2\| x\| \\ &\geq (1 + \varepsilon)^{-1} > 1 - \varepsilon \end{aligned}$$

□

Proposition 5. *If $T \in \mathcal{B}_0(\mathcal{H})$ and $\{\lambda_n\}$ is a sequence of distinct elements in $\sigma_p(T)$, then $\lim_{n \rightarrow \infty} \lambda_n = 0$.*

Proof. For each n , choosing a nonzero $x_n \in \ker(T - \lambda_n)$,

$$\mathcal{M}_n = \text{span} \{ x_1, x_2, \dots, x_n \}$$

By preceding lemma, there is a $y_n \in \mathcal{M}_n$ with $\|y_n\| = 1$ s.t.

$$\text{dist}(y_n, \mathcal{M}_{n-1}) > \frac{1}{2}$$

Let $y_n = \sum_{i=1}^n \alpha_i x_i$, thus

$$(T - \lambda_n)y_n = \sum_{i=1}^{n-1} \alpha_i (\lambda_i - \lambda_n) x_i \in \mathcal{M}_{n-1}$$

Therefore, for $n > m$,

$$\begin{aligned} T(\lambda_n^{-1} y_n) - T(\lambda_m^{-1} y_m) &= \lambda_n^{-1} (T - \lambda_n) y_n - \lambda_m^{-1} (T - \lambda_m) y_m + (y_n - y_m) \\ &= y_n - (y_m + \lambda_n^{-1} (T - \lambda_n) y_n + \lambda_m^{-1} (T - \lambda_m) y_m) \end{aligned}$$

Since the part in the bracketed is in \mathcal{M}_{n-1} ,

$$\|T(\lambda_n^{-1} y_n) - T(\lambda_m^{-1} y_m)\| \leq \text{dist}(y_n, \mathcal{M}_{n-1}) > \frac{1}{2}$$

Thus $\{\lambda_n^{-1} y_n\}$ has no bounded subset by the fact that T is compact.

$$\|\lambda_n^{-1} y_n\| = |\lambda_n|^{-1} \rightarrow \infty$$

That means $\lim_{n \rightarrow \infty} \lambda_n = 0$.

□

Proposition 6. *If $T \in \mathcal{B}_0(\mathcal{H})$ and a nonzero $\lambda \in \sigma(T)$, then λ is a isolated point in $\sigma(T)$.*

Proof. If there is a sequence $\{\lambda_n\} \subset \sigma(T)$ s.t. $\lambda_n \rightarrow \lambda$, then each λ_n is in either $\sigma_p(T)$ or $\sigma_p(T^*)$. Then if there exists a subsequence $\{\lambda_{n_k}\} \subset \sigma_p(T)$. But by above proposition, $\lim_{k \rightarrow \infty} \lambda_{n_k} = 0$, which is a contradiction. Similarly, if $\{\lambda_{n_k}\} \subset \sigma_p(T^*)$, $\lim_{k \rightarrow \infty} \lambda_{n_k} = 0$, which is also a contradiction. \square

Now, we can see any nonzero point in the spectrum of a compact operator is a eigenvalue.

Lemma 7. *If $T \in \mathcal{B}_0(\mathcal{H})$ and a nonzero $\lambda \in \sigma(T)$, then $\ker(T - \lambda)$ is finite dimensional.*

Proof. If there is an infinite orthonormal sequence $\{e_n\}$ in $\ker(T - \lambda)$, then $\{Te_n\}$ has a convergent subsequence $\{Te_{n_k}\}$, thus it is Cauchy.

$$\|Te_{n_k} - Te_{n_j}\|^2 = \|\lambda e_{n_k} - \lambda e_{n_j}\|^2 = 2|\lambda| > 0$$

which is a contradiction. \square

Theorem 8. *If $T \in \mathcal{B}_0(\mathcal{H})$ and a nonzero $\lambda \in \sigma(T)$, then $\lambda \in \sigma_p(T)$ and $\dim \ker(T - \lambda) = \dim \ker(T - \lambda)^* < \infty$ and $\text{ran}(T - \lambda)$ is closed.*

Proof. By the **Proposition 3** in the subsection **2.1.3**, since each $\lambda \in \sigma(T)$ is isolated, we can define

$$E(\lambda) = \int_{\Gamma(\{\lambda\})} (z - T)^{-1} dz$$

where $\Gamma(\{\lambda\})$ is the closed curve enclosed $\{\lambda\}$ and disjoint with other λ in $\sigma(T)$. Then, each $E(\lambda)$ is a projection.

$$\mathcal{H}_\lambda = E(\lambda)\mathcal{H}, \quad T_\lambda = T|_{\mathcal{H}_\lambda}: \mathcal{H}_\lambda \rightarrow \mathcal{H}_\lambda$$

Since $\sigma(T_\lambda) = \{\lambda\}$ and $\lambda \neq 0$, T_λ is invertible. But clearly, T_λ is compact. That means any bounded and closed subset in \mathcal{H}_λ is compact, thus $\dim \mathcal{H}_\lambda < \infty$ by the Riesz Theorem. Therefore, by the result of finite linear algebra, λ is the eigenvalue of T_λ , i.e. $\lambda \in \sigma_p(T)$. By above lemma, $\dim \ker(T - \lambda) < \infty$.

Let $\Delta = \sigma(T) \setminus \{\lambda\}$, $\mathcal{H}_\Delta = E(\{\Delta\})\mathcal{H}$ and $T_\Delta = T|_{\mathcal{H}_\Delta}$. Then since $\lambda \notin \Delta$,

$$\text{ran}(T_\Delta - \lambda) = \mathcal{H}_\Delta$$

Therefore,

$$\begin{aligned} \text{ran}(T - \lambda) &= (T - \lambda)\mathcal{H}_\lambda + (T - \lambda)\mathcal{H}_\Delta \\ &= \text{ran}(T_\lambda - \lambda) + \mathcal{H}_\Delta \end{aligned}$$

Moreover, since $\dim \mathcal{H}_\lambda < \infty$, $\text{ran}(T_\lambda - \lambda)$ is closed. Thus $\text{ran}(T - \lambda)$ is closed.

Finally, note that

$$\begin{aligned}\mathcal{H} / \text{ran}(T - \lambda) &= (\mathcal{H}_\Delta + \mathcal{H}_\lambda) / (\text{ran}(T_\lambda - \lambda) + \mathcal{H}_\Delta) \\ &\cong \mathcal{H}_\lambda / \text{ran}(T_\lambda - \lambda)\end{aligned}$$

Since $\dim \mathcal{H}_\lambda < \infty$,

$$\begin{aligned}\dim(\mathcal{H} / \text{ran}(T - \lambda)) &= \dim \mathcal{H}_\lambda - \dim \text{ran}(T_\lambda - \lambda) = \dim \ker(T_\lambda - \lambda) \\ &= \dim \ker(T - \lambda) < \infty\end{aligned}$$

And since

$$(\mathcal{H} / \text{ran}(T - \lambda))^* \cong \text{ran}(T - \lambda)^\perp = \ker(T - \lambda)^*$$

$$\dim \ker(T - \lambda) = \dim \ker(T - \lambda)^*.$$

□

Combining all of above results, we can get an explicit structure of the spectrum of compact operator.

Theorem 9 (Riesz). *If \mathcal{H} is a infinite dimensional Hilbert space and $T \in \mathcal{B}_0(\mathcal{H})$, then one and only one of the following possibilities occurs.*

- 1) $\sigma(T) = \{0\}$.
- 2) $\sigma(T) = \{0, \lambda_1, \dots, \lambda_n\}$, where each $\lambda_k \neq 0$, and each λ_k is the eigenvalue of T with $\dim \ker(T - \lambda_k) < \infty$.
- 3) $\sigma(T) = \{0, \lambda_1, \lambda_2, \dots\}$, where each $\lambda_k \neq 0$, and each λ_k is the eigenvalue of T with $\dim \ker(T - \lambda_k) < \infty$, and moreover, $\lim_{n \rightarrow \infty} \lambda_n = 0$.

Because of its discrete spectrum, we can using the spectral measure to decompose the compact normal operator.

Theorem 10 (Spectral Theorem). *If T is a compact normal operator on a Hilbert space \mathcal{H} , then T has at most countable eigenvalues $\{\lambda_n\}$, and there are corresponding projections $P_n: \mathcal{H} \rightarrow \ker(T - \lambda_n)$ with $P_n P_m = P_m P_n = 0$ s.t.*

$$T = \sum_{n=1}^{\infty} \lambda_n P_n$$

where $\{\lambda_n\}$ are all distinctive eigenvalues and the convergence is with respect to the norm topology.

Proof. Let E be the spectral measure of T , then we see

$$T = \int_{\sigma(T)} z dE$$

If $\dim \mathcal{H} = \infty$, since $\sigma(N)$ is consisted of at most countable eigenvalues $\{\lambda_n\}_{n=1}^{\infty}$ and 0, then we can set $P_n = E(\{\lambda_n\})$ for $n = 1, 2, \dots$, which is well-defined by the **Proposition 2** in the subsection **3.1.2**, above integral can be

$$\begin{aligned} N &= \lim_{n \rightarrow \infty} \left(\int_{B_{\frac{1}{n}}(0)} z dz + \sum_{k=1}^n \lambda_k P_k \right) \\ &= \sum_{n=1}^{\infty} \lambda_n P_n \end{aligned}$$

And moreover, P_n is the projection from \mathcal{H} onto $\ker(T - \lambda_n)$.

If $\dim \mathcal{H} < \infty$, above theorem is clearly true. \square

Also, we can see the Functional Calculus of the compact normal operators. Firstly, Taking same notation as above theorem and let $P_0 = 1 - \sum_{n=1}^{\infty} P_n$, in fact, P_0 is the projection from \mathcal{H} to $\ker T$. By the Functional Calculus of normal operators, we know

$$W^*(T) = \{ \phi(T) : \phi \in L^{\infty}(\sigma(T)) \}$$

Since $\sigma(T)$ is discrete, $L^{\infty}(\sigma(T)) = l^{\infty}(\mathbb{C})$, and for any $(a_n)_{n=0}^{\infty} \in l^{\infty}(\mathbb{C})$,

$$(a_n)(T) = a_0 P_0 + \sum_{n=1}^{\infty} a_n P_n$$

Thus we can see the separating vector for $W^*(T)$, if e is a separating vector and decomposing e to $\{P_n\}_{n=0}^{\infty}$

$$e = \sum_{n=0}^{\infty} e_n, \quad e_n \in P_n \text{ for each } n$$

Then for $(a_n)_{n=0}^{\infty} \in l^{\infty}(\mathbb{C})$

$$(a_n)(T)e = (a_0 P_0 + \sum_{n=1}^{\infty} a_n P_n) \left(\sum_{n=0}^{\infty} e_n \right) = \sum_{n=0}^{\infty} a_n e_n$$

Since e is separating, $\sum_{n=0}^{\infty} a_n e_n = 0$ for any $(a_n) \in l^{\infty}(\mathbb{C})$ implies $a_n = 0$ for all n . That means that $e_n \neq 0$ for any n . Then the corresponding measure μ is defined as

$$\mu(\{\lambda_n, \lambda_m\}) = \|E(\{\lambda_n, \lambda_m\})e\|^2 = \|e_n + e_m\|^2$$

Now, we can see the multiplicity function of compact normal operator. By similar construction in the **Theorem 3** and **Theorem 8** in the section **3.4**, the multiplicity function is like

$$m_T(\lambda) = \dim \ker(T - \lambda)$$

That means two compact normal operators are equivalent if and only if they have same dimension of all eigenspaces.

Finally,, we can prove $\mathcal{B}_0(\mathcal{H})$ is the unique closed ideal in $\mathcal{B}(\mathcal{H})$ if \mathcal{H} is separable.

Proposition 11. *If N is a normal operator in $\mathcal{B}(\mathcal{H})$ with the spectre measure E , then N is compact if and only if for any $\varepsilon > 0$,*

$$\dim E(\{z \in \mathbb{C} : |z| > \varepsilon\}) < \infty$$

Proof. Let $\Delta_\varepsilon = \{z \in \mathbb{C} : |z| > \varepsilon\}$ and $E_\varepsilon = E(\Delta_\varepsilon)$. Assume that for any given $\varepsilon > 0$, $\dim E_\varepsilon < \infty$, then

$$\begin{aligned} N - NE_\varepsilon &= \int z dE - \int z \chi_{\Delta_\varepsilon} dE \\ &= \int z \chi_{\mathbb{C} \setminus \Delta_\varepsilon} dE = \phi(N) \end{aligned}$$

where $\phi(z) = z \chi_{\mathbb{C} \setminus \Delta_\varepsilon}(z)$. Therefore,

$$\|N - NE_\varepsilon\| = \sup \{|z| : \mathbb{C} \setminus \Delta_\varepsilon\} < \varepsilon$$

Thus $N \in \mathcal{B}_0(\mathcal{H})$.

Conversely, if N is compact, then for any $\varepsilon > 0$, put $\phi(z) = z^{-1} \chi_{\Delta_\varepsilon}$, then

$$N\phi(N) = \int \chi_{\Delta_\varepsilon} dE = E_\varepsilon$$

Since E_ε is a compact projection, i.e. $\text{ran } E_\varepsilon$ is closed, by the Riesz Theorem, $\dim \text{ran } E_\varepsilon < \infty$. \square

Theorem 12. *If \mathcal{H} is a separable Hilbert space and \mathfrak{I} is an ideal of $\mathcal{B}(\mathcal{H})$ that contains a non-compact operator, then $\mathfrak{I} = \mathcal{B}(\mathcal{H})$.*

Proof. Let $T \in \mathfrak{I} \setminus \mathcal{B}_0(\mathcal{H})$ then

$$T^*T = \int_{\sigma(T^*T)} t dE(t)$$

By above theorem, there is an $\varepsilon > 0$ s.t. $P = E(\varepsilon, \infty)$ has infinite rank.

$$P = \left(\int t^{-1} \chi_{(\varepsilon, \infty)}(t) dE(t) \right) T^*T \in \mathfrak{I}$$

Since \mathcal{H} is separable, $\dim P\mathcal{H} = \dim \mathcal{H} = \aleph_0$, there is a unitary U from \mathcal{H} to $P\mathcal{H}$. Therefore, $1 = U^*PU \in \mathfrak{I}$. $\mathfrak{I} = \mathcal{B}(\mathcal{H})$. \square

Proposition 13. *If \mathfrak{I} is a closed ideal of $\mathcal{B}(\mathcal{H})$, the $\mathcal{B}_0(\mathcal{H}) \subset \mathfrak{I}$ or $\mathfrak{I} = \{0\}$.*

Proof. Since \mathfrak{I} is closed, \mathfrak{I} is self-adjoint by the **Proposition 1** in the subsection **2.2.5**. Then by the **Theorem 3** in the subsection **2.2.5**, we know $\mathfrak{I} = \mathfrak{I} \cap \mathfrak{I}^*$ is a hereditary subalgebra. Therefore, if \mathfrak{I} is nonzero and $T \in \mathfrak{I}$ is nonzero, then there is a finite rank projection P s.t.

$$0 \leq P \leq T^*T$$

Thus, $P \in \mathfrak{I}$. Moreover, any finite rank projection is in \mathfrak{I} . Since \mathfrak{I} is norm closed, $\mathcal{B}_0(\mathcal{H}) \subset \mathfrak{I}$. \square

Combining above propositions and theorems, we can get the final result.

Corollary 14. *If \mathcal{H} is a separable Hilbert space, then the only nontrivial closed ideal of $\mathcal{B}(\mathcal{H})$ is $\mathcal{B}_0(\mathcal{H})$.*

5.2 Compact Operator Algebras

For a normal operator, it can be decomposed as the direct sum of *-cyclic normal operators. Therefore, we want to research the similar property of compact operators. Like researching the normal operator, we firstly begin with the representation of compact operators. Let \mathfrak{A} denote a C^* -subalgebra in $\mathcal{B}_0(\mathcal{H})$.

5.2.1 Minimal Projections

Definition 1. If \mathfrak{A} is a C^* -subalgebra in $\mathcal{B}_0(\mathcal{H})$

- 1) A projection E in \mathfrak{A} is minimal if $E \neq 0$ and there are no nonzero projection P in \mathfrak{A} s.t. $P < E$.
- 2) \mathfrak{A} is called irreducible if \mathfrak{A} has no proper reducing subspaces.

Proposition 2. *Let E be a projection in \mathfrak{A} .*

- 1) *E is minimal if and only if*

$$E\mathfrak{A}E = \{ \lambda E : \lambda \in \mathbb{C} \}$$

- 2) *Every projection in \mathfrak{A} is the direct sum of a finite number of pairwise orthogonal minimal projections in \mathfrak{A} .*

Proof. For 1), clearly, $E\mathfrak{A}E$ is a C^* -subalgebra of $\mathcal{B}_0(\mathcal{H})$. Thus for any self-adjoint operator $A \in E\mathfrak{A}E$, each nonzero spectrum point λ of A is isolated. That means $\chi_{\{\lambda\}}$ is continuous, thus $E(\{\lambda\})$ is a projection in $E\mathfrak{A}E$. However, the minimality of E implies that $E(\{\lambda\}) = E$. Therefore, by the fact that $E\mathfrak{A}E$ can be generated by all self-adjoint operators in $E\mathfrak{A}E$, $E\mathfrak{A}E = \{\lambda E : \lambda \in \mathbb{C}\}$. The converse is trivial.

For 2), since each projection in $P \mathfrak{A}$ is compact, P has finite rank. Thus P can be direct sum of a finite number of pairwise orthogonal minimal projections. \square

Using the minimal projection, we can get an important trait the irreducible compact subalgebra has.

Theorem 3. *If \mathfrak{A} is an irreducible C^* -subalgebra of $\mathcal{B}_0(\mathcal{H})$, then $\mathfrak{A} = \mathcal{B}_0(\mathcal{H})$.*

Proof. Let E be a minimal projection in \mathfrak{A} .

Claim: The rank of E is one.

If $\text{rank } E > 1$, then there are nonzero $g, h \in E$ with $g \perp h$. Let $A \in \mathfrak{A}$ and $\lambda \in \mathbb{C}$ s.t. $EAE = \lambda A$. Then

$$\langle Ag, h \rangle = \langle EAEg, h \rangle = \lambda \langle g, h \rangle = 0$$

Thus, $h \perp \overline{\mathfrak{A}g}$. But since \mathfrak{A} is irreducible, $\overline{\mathfrak{A}g} = \mathcal{H}$. So $h = 0$ contradicted to the assumption.

Claim: If P is any rank one projection, then $P \in \mathfrak{A}$.

Let $e \in E$ and $p \in P$. By the irreducibility of \mathfrak{A} , $\mathcal{H} = \overline{\mathfrak{A}e}$. Then for any given $\varepsilon > 0$,

$$\|Ae - p\| < \varepsilon, \quad \text{for some } A \in \mathfrak{A}$$

Let $P' = AE$, then $P' \in \mathfrak{A}$. And moreover, for any $h \in \mathcal{H}$

$$\begin{aligned} \|(P - P')h\| &= \|\langle h, p \rangle p - \langle h, Ae \rangle Ae\| \\ &= \|\langle h, p - Ae \rangle Ae + \langle h, p \rangle (p - Ae)\| \\ &\leq \|h\| \|Ae\| \varepsilon + \|h\| \|p\| \varepsilon \end{aligned}$$

Thus by the fact that \mathfrak{A} is norm closed, $P \in \mathfrak{A}$.

Therefore, any finite rank projection in \mathfrak{A} , i.e. $\mathfrak{A} = \mathcal{B}_0(\mathcal{H})$. \square

Corollary 4. *For any Hilbert space \mathcal{H} , $\mathcal{B}_0(\mathcal{H})$ is simple.*

Proof. Let \mathfrak{I} be an closed ideal in $\mathcal{B}_0(\mathcal{H})$. Then $\overline{\mathfrak{I}e}$ for any nonzero e is a reducing space for $\mathcal{B}_0(\mathcal{H})$. Thus $\overline{\mathfrak{I}e} = \mathcal{H}$ since $\mathcal{B}_0(\mathcal{H})$ is irreducible. That means \mathfrak{I} is irreducible. Thus by above theorem, $\mathfrak{I} = \mathcal{B}_0(\mathcal{H})$. \square

Remark. In fact, the converse of the result in above proof is also true. \mathfrak{B} is a irreducible subalgebra contained in $\mathcal{B}(\mathcal{H})$ if and only if $\overline{\mathfrak{B}e} = \mathcal{H}$ for any $e \in \mathcal{H}$. And because of this, we can get the genera version of the corollary that if \mathfrak{B} is a irreducible subalgebra contained in $\mathcal{B}(\mathcal{H})$, then any ideal of it is irreducible.

Corollary 5. *Let \mathcal{H} be a Hilbert space. If \mathfrak{B} is a irreducible C^* -subalgebra of $\mathcal{B}(\mathcal{H})$ and $\mathfrak{B} \cap \mathcal{B}_0(\mathcal{H}) \neq 0$, then $\mathcal{B}_0(\mathcal{H}) \subset \mathfrak{B}$.*

Proof. Let $\mathfrak{J} = \mathfrak{B} \cap \mathcal{B}_0(\mathcal{H})$. Then \mathfrak{J} is a closed ideal in \mathfrak{B} . By above corollary, $\mathfrak{J} = \mathcal{B}_0(\mathcal{H})$, i.e. $\mathcal{B}_0(\mathcal{H}) \subset \mathfrak{B}$. \square

Now, the unit element of a compact operator subalgebra can be provided.

Proposition 6. *If E is a minimal projection in \mathfrak{A} , e is any nonzero vector in \mathcal{H} and $\mathcal{H}_e = \overline{\mathfrak{A}e}$, then $\mathfrak{A}|_{\mathcal{H}_e} = \mathcal{B}_0(\mathcal{H}_e)$.*

Proof. Let $\mathfrak{A}_e = \mathfrak{A}|_{\mathcal{H}_e}$. We just need to show \mathfrak{A}_e is irreducible in $\mathcal{B}(\mathcal{H}_e)$. Let P be a projection in $\mathcal{B}(\mathcal{H}_e)$ commuting with \mathfrak{A}_e . Put

$$T = P - \langle Pe, e \rangle 1$$

Then $T \in \mathfrak{A}'_e$, for any $A, B \in \mathfrak{A}$,

$$\begin{aligned} \langle TAe, Be \rangle &= \langle TB^*AEe, Ee \rangle \\ &= \langle TEB^*AEe, e \rangle \end{aligned}$$

Since E is minimal, $EB^*AE = \lambda E$ for some $\lambda \in \mathbb{C}$. Therefore,

$$\langle TAe, Be \rangle = \lambda \langle Te, e \rangle = \lambda(\langle Pe, e \rangle - \langle Pe, e \rangle) = 0$$

Since A, B are arbitrary and $\mathcal{H}_e = \overline{\mathfrak{A}e}$, $T = 0$, that means $P = \langle Pe, e \rangle 1$. Thus, $P = 0$ or 1 . \square

5.2.2 Representations of Compact Operator Algebras

For the representation of a compact operator subalgebra, it can be decomposed as the direct sum of irreducible parts.

Theorem 1. *If (ρ, \mathcal{K}) is a non-degenerate representation of \mathfrak{A} , then there are irreducible representations $\{\rho_i\}_{i \in I}$ s.t. each ρ_i is equivalent to a subrepresentation of the identity representation $i: \mathfrak{A} \rightarrow \mathcal{B}(\mathcal{H})$.*

Proof. Claim: There is a minimal projection $E \in \mathfrak{A}$ s.t. $\rho(E) \neq 0$.

In fact, there is a self-adjoint operator $A \in \mathfrak{A}$ s.t. $\rho(A) \neq 0$. Otherwise, ρ is the zero representation. Then there is a spectral projection F for A with $\rho(F) \neq 0$. Then since F is the direct sum of finite numbers minimal projections, that means ρ must not vanish on at least one of them.

For this E , choosing a unit vector $e_0 \in \rho(E)$ and define $\mathcal{K}_0 = \overline{\rho(\mathfrak{A})e_0}$ and $\rho_0(A) = \rho(A)|_{\mathcal{K}_0}$. Then

Claim: (ρ_0, \mathcal{K}_0) is equivalent to a subrepresentation of the identity representation of \mathfrak{A} .

Define the unitary U from \mathcal{K}_0 to $\mathcal{H}_0 = \overline{\mathfrak{A}e}$, where e is the unit vector in E .

$$\begin{aligned} U: \quad \mathcal{K}_0 &\longrightarrow \mathcal{H}_0 \\ \rho(A)e_0 &\longmapsto Ae \end{aligned}$$

Since E is minimal, there exists a unique λ s.t. $EA^*AE = \lambda E$ for $A \in \mathfrak{A}$.

$$\begin{aligned}\|\rho(A)e_0\|^2 &= \|\rho(AE)e_0\|^2 = \langle \rho(EA^*AE)e_0, e_0 \rangle \\ &= \lambda \langle e_0, e_0 \rangle = \lambda \langle e, e \rangle \\ &= \langle EA^*AEe_0, e_0 \rangle \\ &= \|Ae\|^2\end{aligned}$$

Therefore, U is surjective isometry, and can extend to \mathcal{H}_0 . And moreover,

$$U\rho(A)\rho(B)e_0 = ABe \Rightarrow U\rho(A)U^*Be = A|_{\mathcal{H}_0}Be$$

Thus, $U\rho_0(A)U^* = A|_{\mathcal{H}_0}$. By above proposition, ρ_0 is irreducible.

Then by Zorn's Lemma, we can get a maximal family of $(\rho_i, \mathcal{H}_i)_{i \in I}$ as above construction. And by the maximality and the non-degenerality,

$$\rho = \oplus_{i \in I} \rho_i, \quad \mathcal{K} = \oplus_{i \in I} \mathcal{H}_i \quad \square$$

Remark. If ρ is degenerate, then let $\mathcal{H}_0 = \overline{\rho(\mathfrak{A})\mathcal{H}}$. The representation $(\rho_0(A) = \rho(A)|_{\mathcal{H}_0}, \mathcal{H}_0)$ is a non-degenerate representation, thus above theorem can be applied to it. Moreover, $\rho = \rho_0 \oplus 0$.

Now, we can use above theorem to get the structure of any finite dimensional C^* -algebra.

Corollary 2. *For any finite dimensional C^* -algebra \mathfrak{A} and let $M(n)$ be the set of $n \times n$ matrices with complex entries acting on the inner product space \mathbb{C}^n , there are $n_1, \dots, n_p \in \mathbb{N}$, s.t.*

$$\mathfrak{A} \cong M(n_1) \oplus M(n_2) \oplus \dots \oplus M(n_p)$$

Also, there are some other corollaries of above theorem.

Corollary 3. 1) *If (ρ, \mathcal{K}) is an irreducible representation of \mathfrak{A} , then $\rho(\mathfrak{A}) = \mathcal{B}_0(\mathcal{K})$.*

2) *If \mathcal{H} and \mathcal{K} are Hilbert spaces and $\rho: \mathcal{B}_0(\mathcal{H}) \rightarrow \mathcal{B}_0(K)$ is an $*$ -isomorphism, then there is a unitary $U: \mathcal{H} \rightarrow \mathcal{K}$, s.t.*

$$\rho(T) = UTU^*, \text{ for any } T \in \mathcal{B}_0(\mathcal{H})$$

3) *If \mathcal{H} and \mathcal{K} are Hilbert spaces and $\rho: \mathcal{B}(\mathcal{H}) \rightarrow \mathcal{B}(K)$ is an $*$ -isomorphism, then there is a unitary $U: \mathcal{H} \rightarrow \mathcal{K}$, s.t.*

$$\rho(A) = UAU^*, \text{ for any } A \in \mathcal{B}(\mathcal{H})$$

Proof. 1) is the direct result from above theorem. And 2) can be obtained by 1).

For 3), since $\rho(\mathcal{B}_0(\mathcal{H}))$ is an ideal of \mathcal{H} and \mathcal{H} is irreducible, $\rho(\mathcal{B}_0(\mathcal{H}))$ is irreducible by above corollary. Thus $\rho(\mathcal{B}_0(\mathcal{H})) = \mathcal{B}_0(\mathcal{H})$. Then by 2), there is a unitary $U: \mathcal{H} \rightarrow \mathcal{H}$, s.t.

$$\rho(A) = UAU^*, \text{ for any } A \in \mathcal{B}_0(\mathcal{H})$$

Let $\{E_i\}$ be the approximatel identity for $\mathcal{B}_0(\mathcal{H})$ consisting with all finite projections. For any $A \leq 0$ in \mathcal{H} , then it can see

$$A_i = A^{\frac{1}{2}} E_i A^{\frac{1}{2}} \rightarrow A \text{ in } SOT$$

Therefore, $\rho(A_i) = UA_iU^* \rightarrow UAU^* = T$ in SOT . Since $A_i \leq A$, $\rho(A_i) \leq \rho(A)$ and thus $T \leq \rho(A)$ by the fact that $\{A_i\}$ is increasing. Conversely, $\rho(A_i) \leq T$ implies that $A_i \leq \rho^{-1}(T)$. Thus $A \leq \rho^{-1}(T)$, thus $\rho(A) \leq T$. \square

Here is an important example, which can be used to construct the AF algebra.

Example 4. Let $M(n)$ be the set of $n \times n$ matrices with complex entries acting on the inner product space \mathbb{C}^n . If

$$\rho: M(m) \longrightarrow M(n)$$

is a $*$ -homomorphism, then there is an integer k s.t. $km \leq n$ and a unitary U in $M(n)$ and

$$\rho(x) = U(\underbrace{x \oplus x \oplus \cdots \oplus x}_k \oplus 0)U^*$$

In general, define $\mathbf{m} = (m_1, m_2, \dots, m_p)$, where $m_i \in \mathbb{N}$, and

$$M(\mathbf{m}) = M(m_1) \oplus M(m_2) \oplus \cdots M(m_p)$$

Then for $\mathbf{m} = (m_1, m_2, \dots, m_p)$ and $\mathbf{n} = (n_1, n_2, \dots, n_q)$, if $\rho: M(\mathbf{m}) \rightarrow M(\mathbf{n})$ is a $*$ -homomorphism, then there this a $q \times p$ matirx $[k_{ij}]$ with integer entries s.t.

$$\rho(x_1, \dots, x_p) = U_1(\underbrace{x_1 \oplus \cdots \oplus x_1}_{k_{11}} \oplus 0)U_1^* \oplus \cdots \oplus U_q(\underbrace{x_p \oplus \cdots \oplus x_p}_{k_{qp}} \oplus 0)U_q^*$$

By above corollary, we see for any finite dimensional C^* -algebras \mathfrak{A} , there exists a $\mathbf{n} = (n_1, n_2, \dots, n_q)$ s.t. $\mathfrak{A} \cong M(\mathbf{n})$. Then we can get the form of the $*$ -homomorphism between any two finite dimensional C^* -algebras.

If ρ is a $*$ -homomorphism between two finite dimensional C^* -algebras $\mathfrak{A} \cong M(\mathbf{m})$ and $\mathfrak{B} \cong M(\mathbf{n})$, where $\mathbf{m} = (m_1, m_2, \dots, m_p)$ and $\mathbf{n} = (n_1, n_2, \dots, n_q)$, then ρ is determined up to a $q \times p$ matirx $[k_{ij}]$ with integer entries. In fact, we can use it to construct the AF algebras.

5.2.3 Decompositions of Compact Operator Algebras

By above theorem and the Zorn's Lemma, for the non-degenerate \mathfrak{A} , there is a maximal family $\{E_i\}_{i \in I}$ of pairwise orthogonal minimal projections in \mathfrak{A} and let $\mathcal{H}_i = \overline{\mathfrak{A}E_i\mathcal{H}}$ s.t.

$$\mathcal{H} = \oplus_{i \in I} \mathcal{H}_i, \quad \mathfrak{A}|_{\mathcal{H}_i} = \mathcal{B}_0(\mathcal{H}_i)$$

But we cannot just using this as the decomposition of the compact operator algebra \mathfrak{A} , since there are some "equivalent" relationships of these \mathcal{H}_i . Firstly, we say that \mathcal{H}_i does not dependent on \mathcal{H}_j , if there is a $A \in \mathfrak{A}$ s.t. $\mathfrak{A}|_{\mathcal{H}_i} = 0$ and $\mathfrak{A}|_{\mathcal{H}_j} \neq 0$.

Proposition 1. \mathcal{H}_i does not dependent on \mathcal{H}_j if and only if

$$A|_{\mathcal{H}_i \oplus \mathcal{H}_j} = \mathcal{B}_0(\mathcal{H}_i) \oplus \mathcal{B}_0(\mathcal{H}_j)$$

Proof. Assume that \mathcal{H}_i does not dependent on \mathcal{H}_j . Let

$$\mathfrak{I} = \{ A \in \mathfrak{A} : A|_{\mathcal{H}_j} = 0 \}$$

Since \mathcal{H}_j is reducing, \mathfrak{I} is a closed ideal. Moreover, $\mathfrak{I}|_{\mathcal{H}_i}$ is a closed ideal of $\mathfrak{A}|_{\mathcal{H}_i} = \mathcal{B}_0(\mathcal{H}_i)$. By the assumption, $\mathfrak{I}|_{\mathcal{H}_i}$ is not zero, thus $\mathfrak{I}|_{\mathcal{H}_i} = \mathcal{B}_0(\mathcal{H}_i)$. Let K_i and K_j be arbitrary compact operators on \mathcal{H}_i and \mathcal{H}_j . And $A \in \mathfrak{A}$ s.t. $A|_{\mathcal{H}_j} = K_j$. Since $\mathfrak{I}|_{\mathcal{H}_i} = \mathcal{B}_0(\mathcal{H}_i)$, there exists $B \in \mathfrak{I}$ s.t. $B|_{\mathcal{H}_i} = K_i - A|_{\mathcal{H}_i}$. Thus,

$$(A + B)|_{\mathcal{H}_i} = K_i, \quad \& \quad (A + B)|_{\mathcal{H}_j} = K_j$$

The converse is trivial. □

Remark. By this proposition, if \mathcal{H}_i does not dependent with \mathcal{H}_j , \mathcal{H}_j does not \mathcal{H}_i . \mathcal{H}_i and \mathcal{H}_j are independent if they are not dependent.

Now we want to research the same properties of the dependent subspaces have, and then we can divide all $\mathcal{H}_{i \in I}$ into dependency class.

Proposition 2. If \mathcal{H}_i and \mathcal{H}_j are dependent, then

$$\|A|_{\mathcal{H}_i}\| = \|A|_{\mathcal{H}_j}\|, \quad \text{for any } A \in \mathfrak{A}$$

Moreover, two dependent spaces are isomorphic, and each dependency class is finite.

Proof. By the hypothesis, if $A \in \mathfrak{A}$ s.t. $A|_{\mathcal{H}_i} = 0$, then $A|_{\mathcal{H}_j} = 0$. Then we can define the map

$$\rho: \mathcal{B}_0(\mathcal{H}_j) \longrightarrow \mathcal{B}_0(\mathcal{H}_i)$$

For any $K \in \mathcal{B}_0(\mathcal{H}_j)$, there is a $A \in \mathfrak{A}$ s.t. $A|_{\mathcal{H}_j} = K$, then we define that $\rho(K) = A|_{\mathcal{H}_i}$.

Claim: ρ is well-defined.

If $B \in \mathfrak{A}$ s.t. $B|_{\mathcal{H}_j} = K$, then $(A - B)|_{\mathcal{H}_j} = 0$. By the dependency of \mathcal{H}_j and \mathcal{H}_i , $(A - B)|_{\mathcal{H}_i} = 0$, i.e. $A|_{\mathcal{H}_i} = B|_{\mathcal{H}_i}$.

ρ is a $*$ -isomorphism.

If $K \in \mathcal{B}_0(\mathcal{H}_j)$ s.t. $\rho(K) = 0$, then there is a $A \in \mathfrak{A}$ s.t. $A|_{\mathcal{H}_j} = K$ and $A|_{\mathcal{H}_i} = 0$. Also, by the dependency, $K = A|_{\mathcal{H}_j} = 0$. Clearly, ρ is a $*$ -homomorphism, then ρ is a $*$ -monomorphism, i.e. ρ is an isometry. And the surjectivity of ρ is similar by the dependency of \mathcal{H}_j and \mathcal{H}_i .

Then by the **Corollary 3** in above subsection, we can find a unitary from \mathcal{H}_j to \mathcal{H}_i . Thus $\dim \mathcal{H}_j = \dim \mathcal{H}_i$. And if there are infinite element in a dependency class, by the compactness of A , $\|A|_{\mathcal{H}_i}\| \rightarrow 0$, which is a contradiction. \square

Theorem 3. *If \mathfrak{A} is a C^* -subalgebra of compact operators, then*

$$\mathcal{H} \cong \mathcal{H}_0 \oplus \oplus_{d \in D} \mathcal{H}_d^{(k_d)}$$

and

$$\mathfrak{A} \cong \{ 0 \oplus \oplus_{d \in D} K_d^{(k_d)} : K_d \in \mathcal{B}_0(\mathcal{H}_d) \}$$

Proof. Firstly, by above mention, if \mathfrak{A} is non-degenerate, there are $\{\mathcal{H}_i\}_{i \in I}$ s.t.

$$\mathcal{H} = \oplus_{i \in I} \mathcal{H}_i, \quad \mathfrak{A}|_{\mathcal{H}_i} = \mathcal{B}_0(\mathcal{H}_i)$$

But by above theorem, we can divide $\{\mathcal{H}_i\}$ into the dependency classes, and re-labeled as

$$\{ \underbrace{\mathcal{H}_1, \dots, \mathcal{H}_1}_{k_1}, \underbrace{\mathcal{H}_2, \dots, \mathcal{H}_2}_{k_2}, \dots \}$$

Thus, for the non-degenerate \mathfrak{A} ,

$$\mathcal{H} \cong \oplus_{d \in D} \mathcal{H}_d^{(k_d)}, \quad \mathfrak{A} \cong \{ \oplus_{d \in D} K_d^{(k_d)} : K_d \in \mathcal{B}_0(\mathcal{H}_d) \}$$

If \mathfrak{A} is degenerate, let $\mathcal{H}_0 = (\overline{\mathfrak{A}\mathcal{H}})^\perp$, then this theorem holds. \square

For a non-degenerate \mathfrak{A} , let $\hat{\mathfrak{A}}$ denote the set of all equivalent classes of irreducible representations of \mathfrak{A} . For any $\zeta \in \hat{\mathfrak{A}}$, let $\rho_\zeta \in \zeta$. By the theorem, for any representation ρ of \mathfrak{A} , there are irreducible representations $\{\rho_i\}$ s.t. $\rho = \oplus_i \rho_i$, then we define the mutiplicity function of ρ ,

$$m_\rho(\zeta) = \#\{i : \rho_i \in \zeta\}$$

Now, we can easily apply above theorem to the representation.

Theorem 4. *If ρ is a representation of \mathfrak{A} and m_ρ is the multiplicity function, then*

$$\rho \cong \oplus \{ \rho_\zeta^{(m_\rho(\zeta))} : \zeta \in \hat{\mathfrak{A}} \}$$

Moreover, any two representations of \mathfrak{A} are equivalent if and only if they have same multiplicity functions.

5.3 Trace Class and Ultraweak Topology

There some interesting subalgebras in the compact operator algebras, and some of them can provide extra topologies on operator algebras, which will be helpful in the research of general von Neumann algebras.

5.3.1 Trace Class and Hilbert-Schmit Operators

Definition 1. Let $T \in \mathcal{B}(\mathcal{H})$. If there is a orthonormal basis \mathcal{E} of \mathcal{H} s.t.

$$\sum_{e \in \mathcal{E}} \langle |T| e, e \rangle < \infty$$

then T is called trace class. Let $\mathcal{B}_1(\mathcal{H})$ denote the set of all trace classes operators.

In this definition, the orthonormal basis just need to exist. But how can that condition garantees for any orthonormal basis of \mathcal{H} the condition is always satisfied?

Proposition 2. If \mathcal{E} and \mathcal{F} are two orthonormal bases for \mathcal{H} , then for any $T \in \mathcal{B}(\mathcal{H})$,

$$\sum_{e \in \mathcal{E}} \|Te\|^2 = \sum_{f \in \mathcal{F}} \|T^*f\|^2 = \sum_{e \in \mathcal{E}} \sum_{f \in \mathcal{F}} |\langle Te, f \rangle|^2$$

Proof. By the Parseval's Identity,

$$\|Te\|^2 = \sum_{f \in \mathcal{F}} |\langle Te, f \rangle|^2$$

Then we get above identity. □

Corollary 3. The sum

$$\sum_{e \in \mathcal{E}} \langle |T| e, e \rangle$$

is independent with the choise of the \mathcal{E} .

Therefore, by this corollary, above definition is valid. Moreover, we can define one more norm on $\mathcal{B}_1(\mathcal{H})$, for $T \in \mathcal{B}_1(\mathcal{H})$ and some orthonormal basis \mathcal{E} ,

$$\|T\|_1 = \sum_{e \in \mathcal{E}} \langle |T| e, e \rangle$$

Then by using this norm, called the trace norm, on $\mathcal{B}_1(\mathcal{H})$, the $\mathcal{B}_1(\mathcal{H})$ can be a Banach space.

Definition 4. $T \in \mathcal{B}(\mathcal{H})$ is called a Hilbert-Schmidt operator if $|T|^2$ is trace class. Let $\mathcal{B}_2(\mathcal{H})$ denote the set of all Hilbert-Schmidt operators.

Also, we can define the norm on $\mathcal{B}_2(\mathcal{H})$, for any orthonormal basis \mathcal{E} in \mathcal{H} ,

$$\|T\|_2 = \left(\sum_{e \in \mathcal{E}} \langle |T|^2 e, e \rangle \right)^{\frac{1}{2}} = \left(\sum_{e \in \mathcal{E}} \| |T| e \|^2 \right)^{\frac{1}{2}} = \left(\sum_{e \in \mathcal{E}} \|Te\|^2 \right)^{\frac{1}{2}} = \left\| |T|^2 \right\|_1^{\frac{1}{2}}$$

And $\mathcal{B}_2(\mathcal{H})$ can be also a Banach space.

In order to research the trace class, it is convinient to get some properties of the Hilbert-Schmidt operators.

Proposition 5. *Let $T \in \mathcal{B}_2(\mathcal{H})$.*

- 1) $\|T\|_2 = \left(\sum_{e \in \mathcal{E}} \|Te\|^2 \right)^{\frac{1}{2}}$.
- 2) $\|T^*\|_2 = \|T\|_2$.
- 3) $\|T\| \leq \|T\|_2$.
- 4) *If $A \in \mathcal{B}(\mathcal{H})$, then $AT, TA \in \mathcal{B}_2(\mathcal{H})$ and*

$$\|AT\|_2, \|TA\|_2 \leq \|A\| \|T\|_2$$

- 5) $\mathcal{B}_2(\mathcal{H})$ *is an ideal of $\mathcal{B}(\mathcal{H})$ and $\|\cdot\|_2$ is a norm on $\mathcal{B}_2(\mathcal{H})$.*

Proof. The first three results are trivial. For 4), fix an orthonormal basis \mathcal{E} and an $A \in \mathcal{B}(\mathcal{H})$, for a $e \in \mathcal{E}$,

$$\|ATe\|^2 \leq \|A\|^2 \|Te\|^2 = \|A\|^2 \| |T| e \|^2$$

Therefore, $\|AT\|_2 \leq \|A\| \|T\|_2$.

For 5), we just need to prove the addition is continuous and closed. For a fixed normal basis \mathcal{E} and $T, S \in \mathcal{B}_2(\mathcal{H})$, then

$$\{\|Te\| : e \in \mathcal{E}\}, \{\|Se\| : e \in \mathcal{E}\} \in l^2(\mathcal{E})$$

By the triangle inequality for $l^2(\mathcal{E})$,

$$\|S + T\|_2 = \left(\sum_E (\|Te\| + \|Se\|)^2 \right)^{\frac{1}{2}} \leq \|T\|_2 + \|S\|_2 \quad \square$$

Remark. Therefore, $\mathcal{B}_2(\mathcal{H})$ with the fixed involution and the norm $\|\cdot\|_2$ is a C^* -algebra.

By using these results, we can see the Hilbert-Schmidt operator is compact.

Corollary 6. *If $T \in \mathcal{B}_2(\mathcal{H})$ and $\varepsilon > 0$, there is a $A \in \mathcal{B}_{00}(\mathcal{H})$ s.t.*

$$\|T - A\|_2 \leq \varepsilon$$

Consequently, every Hilbert-Schmidt operator is compact.

Proof. For this fixed $\varepsilon > 0$, since $T \in \mathcal{B}_2(\mathcal{H})$, there is a finite set $I \in \mathcal{E}$ s.t.

$$\sum_{\mathcal{E} \setminus I} \|Te\|^2 < \varepsilon^2$$

Let $\hat{E} = \text{span } E$ and $B = T|_{\hat{E}} \in \mathcal{B}_{00}(\mathcal{H})$, then

$$\|T - B\|_2 = \left(\sum_{\mathcal{E} \setminus I} \|Te\|^2 \right)^{\frac{1}{2}} < \varepsilon \quad \square$$

Then, there are similar consequences of the $\mathcal{B}_1(\mathcal{H})$.

Proposition 7. *If $T \in \mathcal{B}(\mathcal{H})$, then the following statements are equivalent.*

- 1) $T \in \mathcal{B}_1(\mathcal{H})$.
- 2) $|T|^{\frac{1}{2}} \in \mathcal{B}_2(\mathcal{H})$.
- 3) T is the product of two Hilbert-Schmidt operators.
- 4) $|T|$ is the product of two Hilbert-Schmidt operators.

Proof. Since $\left\| |T|^{\frac{1}{2}} \right\| e^2, 1)$ implies 2). And by the Polar Decomposition,

$$T = W |T| = (W |T|^{\frac{1}{2}}) |T|^{\frac{1}{2}}$$

thus 2) implies 3). And 3) implies 3), because of the Polar Decomposition. Suppose $|T| = BC$ for $B, C \in \mathcal{B}_2(\mathcal{H})$, then for any orthonormal \mathcal{E} and $e \in \mathcal{E}$

$$\langle |T| e, e \rangle = \langle CeB^*e \rangle \leq \|Ce\| \|B^*e\|$$

Therefore, we have

$$\begin{aligned} \sum_{\mathcal{E}} \langle |T| e, e \rangle &\leq \sum_{\mathcal{E}} \|Ce\| \|B^*e\| \\ &\leq \left(\sum_{\mathcal{E}} \|Ce\|^2 \right)^{\frac{1}{2}} \left(\sum_{\mathcal{E}} \|Be\|^2 \right)^{\frac{1}{2}} \\ &= \|C\|_2 \|B\|_2 \end{aligned}$$

Thus, $T \in \mathcal{B}_1(\mathcal{H})$. \square

Definition 8. If $T \in \mathcal{B}_1(\mathcal{H})$ and \mathcal{E} is any orthonormal basis, the trace of T is defined as

$$\text{tr } T = \sum_{\mathcal{E}} \langle Te, e \rangle$$

It need to check that this definition is well-defined.

Proposition 9. *If $T \in \mathcal{B}_1(\mathcal{H})$ and \mathcal{E} is any orthonormal basis, then $\sum_{\mathcal{E}} |\langle Te, e \rangle| < \infty$, and the sum $\sum_{\mathcal{E}} \langle Te, e \rangle$ is independent with the choice of \mathcal{E} .*

Proof. By above proposition, $T = C^*B$ for $C, B \in \mathcal{B}_2(\mathcal{H})$. Since $\|(C - \lambda B)e\|^2 \geq 0$ for any $\lambda \in \mathbb{C}$,

$$2 \operatorname{Re} \bar{\lambda} \langle Be, Ce \rangle \leq \|Be\|^2 + |\lambda| \|Ce\|^2$$

Choosing a λ s.t. $|\lambda| = 1$ and

$$\bar{\lambda} \langle Be, Ce \rangle = |\langle Be, Ce \rangle|$$

Then we have,

$$|\langle Te, e \rangle| = |\langle Be, Ce \rangle| \leq \frac{1}{2}(\|Be\|^2 + \|Ce\|^2)$$

Thus $\sum_{\mathcal{E}} |\langle Te, e \rangle| \leq \frac{1}{2}(\|B\|_2^2 + \|C\|_2^2)$.
And since

$$\begin{aligned} \operatorname{Re} \langle Te, e \rangle &= \frac{1}{4}(\|(B + C)e\|^2 - \|(B - C)e\|^2) \\ \operatorname{Im} \langle Te, e \rangle &= \frac{1}{4}(\|(iB + C)e\|^2 - \|(iB - C)e\|^2) \end{aligned}$$

then

$$\sum_{\mathcal{E}} \langle Te, e \rangle = \operatorname{Re} \sum_{\mathcal{E}} \langle Te, e \rangle + i \operatorname{Im} \sum_{\mathcal{E}} \langle Te, e \rangle$$

The sum is independent with the choice of the orthonormal basis \mathcal{E} . \square

For convenience, there is another notation to denote the rank one operators.

Definition 10. For two vectors $g, h \in \mathcal{H}$, the rank one operator $g \otimes h$ is defined as

$$g \otimes h(f) = \langle f, h \rangle g$$

There are some elementary properties for this rank one operator.

Proposition 11. *Let $g, h \in \mathcal{H}$.*

- 1) $e \otimes e$ is the projection onto $\mathbb{C}e$.
- 2) $\operatorname{ran} g \otimes h = \mathbb{C}g$ and $\ker g \otimes h = (\mathbb{C}h)^\perp$ for $g \neq 0$.
- 3) $(g \otimes h)^* = h \otimes g$.
- 4) The map $(g, h) \rightarrow g \otimes h$ is a sesquilinear from $\mathcal{H} \times \mathcal{H}$ to $\mathcal{B}(\mathcal{H})$.

5) $\|g \otimes h\| = \|g\| \|h\|$.

6) If $T \in \mathcal{B}(\mathcal{H})$, then $T(g \otimes h) = (Tg) \otimes h$ and $(g \otimes h)T = g \otimes (T^*h)$.

7) Every any finite rank operator can be expressed as

$$\sum_{k=1}^n g_k \otimes h_k$$

for some $g_1, \dots, g_n \in \mathcal{H}$ and $h_1, \dots, h_n \in \mathcal{H}$.

8) Every any finite rank operator can be expressed as

$$\sum_{k=1}^n g_k \otimes h_k$$

for some $g_1, \dots, g_n \in \mathcal{H}$ and $h_1, \dots, h_n \in \mathcal{H}$.

9) For $C = \sum_{k=1}^n g_k \otimes h_k$ and $A \in \mathcal{B}(\mathcal{H})$,

$$\text{tr } AC = \sum_{k=1}^n \langle Ag_k, h_k \rangle$$

There are more interesting properties of the trace class.

Theorem 12. 1) $\mathcal{B}_1(\mathcal{H}) \subset \mathcal{B}_0(\mathcal{H})$ and conversely, if $A \in \mathcal{B}_0(\mathcal{H})$ and $\{\alpha_n\}$ are the eigenvalues of $|A|$, then $A \in \mathcal{B}_1(\mathcal{H})$ if and only if $(\alpha_n) \in l^1$. In this case, $\|A\|_1 = \sum \alpha_n$.

2) $\mathcal{B}_1(\mathcal{H})$ is an ideal of $\mathcal{B}(\mathcal{H})$.

3) $\text{tr} : \mathcal{B}_1(\mathcal{H}) \rightarrow \mathbb{C}$ is a positive non-degenerate linear functional.

4) $\mathcal{B}_{00}(\mathcal{H})$ is a dense subset in $\mathcal{B}_1(\mathcal{H})$ with respect to $\|\cdot\|_1$.

5) If $T \in \mathcal{B}_1(\mathcal{H})$, then for any $A \in \mathcal{B}(\mathcal{H})$,

$$\text{tr } TA = \text{tr } AT, \text{ and } |\text{tr } AT| \leq \|A\| \|T\|_1$$

6) $\|T\|_1 = \|T^*\|_1$ for any $T \in \mathcal{B}_1(\mathcal{H})$.

7) If $T \in \mathcal{B}_1(\mathcal{H})$ and $A \in \mathcal{B}(\mathcal{H})$, then

$$\|TA\|_1, \|AT\|_1 \leq \|A\| \|T\|_1$$

Proof. If $A \in \mathcal{B}_0(\mathcal{H})$, also $|A| \in \mathcal{B}_0(\mathcal{H})$, then there is an orthonormal basis $\{e_n\}$ for \mathcal{H} and the corresponding eigenvalues $\{\alpha_n\} \in l^\infty(\mathbb{C})$ s.t.

$$|A| = \sum_n \alpha_n P_{e_n} = \sum_n \alpha_n e_n \otimes e_n$$

Then we have that

$$\sum_n \langle |A| e_n, e_n \rangle = \sum_n \alpha_n$$

By the Polar Decomposition, $A \in \mathcal{B}_1(\mathcal{H})$ is equivalent to $|A| \in \mathcal{B}_1(\mathcal{H})$. Thus, $\{\alpha_n\} \in l^1(\mathbb{C})$ is equivalent to $A \in \mathcal{B}_1(\mathcal{H})$, and $\|A\|_1 = \sum \alpha_n$.

For 2), similarly as the $\mathcal{B}_2(\mathcal{H})$, we just need to prove that the addition is continuous and closed. Let $A \in \mathcal{B}_1(\mathcal{H})$ and $B \in \mathcal{B}_1(\mathcal{H})$, and by the Polar Decomposition,

$$A = W |A|, \quad B = V |B|, \quad A + B = U |A + B|$$

And $|A + B| = \sum_n \gamma_n e_n \otimes e_n$ for an orthonormal basis $\{e_n\}$. By 1), we just need to check that $\{\gamma_n\} \in l^1(\mathbb{C})$.

$$\begin{aligned} \sum_n \gamma_n &= \sum_n \langle |A + B| e_n, e_n \rangle = \sum_n |\langle A e_n, U e_n \rangle + \langle B e_n, U e_n \rangle| \\ &= \sum_n |\langle |A| e_n, W^* U e_n \rangle + \langle |B| e_n, V^* U e_n \rangle| \\ &= \sum_n \left| \langle |A|^{\frac{1}{2}} e_n, |A|^{\frac{1}{2}} W^* U e_n \rangle + \langle |B|^{\frac{1}{2}} e_n, |B|^{\frac{1}{2}} V^* U e_n \rangle \right| \\ &\leq \sum_n (\| |A|^{\frac{1}{2}} e_n \| \| |A|^{\frac{1}{2}} W^* U e_n \| + \| |B|^{\frac{1}{2}} e_n \| \| |B|^{\frac{1}{2}} V^* U e_n \|) \\ &\leq \left(\sum_n \| |A|^{\frac{1}{2}} e_n \|^2 \right)^{\frac{1}{2}} \left(\sum_n \| |A|^{\frac{1}{2}} W^* U e_n \|^2 \right)^{\frac{1}{2}} \\ &\quad + \left(\sum_n \| |B|^{\frac{1}{2}} e_n \|^2 \right)^{\frac{1}{2}} \left(\sum_n \| |B|^{\frac{1}{2}} V^* U e_n \|^2 \right)^{\frac{1}{2}} \\ &\leq \| |A|^{\frac{1}{2}} \|_2^2 + \| |B|^{\frac{1}{2}} \|_2^2 \\ &= \|A\|_1 + \|B\|_1 \end{aligned}$$

3) is trivial and 4) can be obtained by similar argument as the $\mathcal{B}_2(\mathcal{H})$.

For 5), let $T \in \mathcal{B}_1(\mathcal{H})$ and $T = C^* B$ for some $B, C \in \mathcal{B}_2(\mathcal{H})$, then by above mention,

$$\begin{aligned} \operatorname{Re} \operatorname{tr} C^* B &= \frac{1}{4} (\|(B + C)e\|^2 - \|(B - C)e\|^2) \\ &= \frac{1}{4} (\|(B^* + C^*)e\|^2 - \|(B^* - C^*)e\|^2) \\ &= \operatorname{Re} \operatorname{tr} C B^* \end{aligned}$$

And similarly, $\text{Im tr } C^*B = -\text{Im tr } CB^*$, thus

$$\text{tr } C^*B = \overline{\text{tr } CB^*}$$

Therefore, we have that for any $A \in \mathcal{B}(\mathcal{H})$

$$\text{tr } AT = \text{tr } (AC^*)B = \overline{\text{tr } CA^*B^*} = \text{tr } (C^*B)A = \text{tr } TA$$

Let $T = W|T|$ be the Polar Decomposition. Then by the CBS Inequality,

$$\begin{aligned} |\text{tr } AT| &\leq \sum_{\mathcal{E}} \left\| |T|^{\frac{1}{2}} e \right\| \left\| |T|^{\frac{1}{2}} W^* A^* e \right\| \\ &\leq \left(\sum \left\| |T|^{\frac{1}{2}} e \right\|^2 \right)^{\frac{1}{2}} \left(\sum \left\| |T|^{\frac{1}{2}} W^* A^* e \right\|^2 \right)^{\frac{1}{2}} \\ &\leq \left\| |T|^{\frac{1}{2}} \right\|_2 \left\| |T|^{\frac{1}{2}} W^* A^* \right\|_2 \\ &\leq \left\| |T|^{\frac{1}{2}} \right\|_2^2 \|W^* A^*\| \\ &\leq \|T\|_1 \|A\| \end{aligned}$$

For 6), let $T = W|T| \in \mathcal{B}_1(\mathcal{H})$, then $TT^* = W|T|^2 W^*$. So by the uniqueness of root, $|T^*| = W|T| W^*$. Therefore,

$$\|T^*\|_1 = \text{tr } |T^*| = \text{tr } W|T| W^* = \text{tr } WW^*|T|$$

Since $WW^*|T| = |T|$, $\|T^*\|_1 = \|T\|_1$.

For 7), let $T = W|T| \in \mathcal{B}_1(\mathcal{H})$ and $AT = V|AT|$. So

$$|AT| = S|T|, \text{ where } S = V^*AW$$

Thus $\|S\| \leq \|A\|$, then

$$\|AT\|_1 = \text{tr } |AT| = \text{tr } S|T| \leq \|S\| \|T\|_1 \leq \|A\| \|T\|_1 \quad \square$$

Remark. By this theorem, $(\mathcal{B}_1(\mathcal{H}), \|\cdot\|_1)$ is also a C^* -algebra.

Theorem 13. *If $\{g_n\}$ and $\{h_n\}$ are two square summable sequences, then*

$$T = \sum_n g_n \otimes h_n \in \mathcal{B}_1(\mathcal{H}), \text{ and } \|T\|_1 \leq \sum_n \|g_n\| \|h_n\|$$

Conversely, if $T \in \mathcal{B}_1(\mathcal{H})$, then there are two square summable sequences $\{g_n\}$ and $\{h_n\}$ s.t.

$$\sum_n \|g_n\|^2 = \|T\|_1 = \sum_n \|h_n\|^2, \text{ and } T = \sum_n g_n \otimes h_n$$

Proof. Let $\{e_n\}$ be an orthonormal sequence of \mathcal{H} and $G_n = \sum_{k=1}^n g_k \otimes e_k$. If $n > m$, then

$$\|G_n - G_m\|_2^2 = \sum_{k=m+1}^n \|(G_n - G_m)e_k\|^2 = \sum_{k=m+1}^n |\langle e_k, g_k \rangle|^2 \leq \sum_{k=m+1}^n \|g_k\|^2$$

Since $\{g_n\}$ is square summable, $\{G_n\}$ is Cauchy with respect to the norm $\|\cdot\|_2$. Therefore, there is a $G \in \mathcal{B}_2(\mathcal{H})$ s.t.

$$G = \lim_n G_n = \sum_n g_n \otimes e_n$$

Similarly, there is a $H = \sum_n e_n \otimes h_n$ in $\mathcal{B}_2(\mathcal{H})$. And then by the definition, $T = GH \in \mathcal{B}_1(\mathcal{H})$.

Conversely, if $T \in \mathcal{B}_1(\mathcal{H})$ and $T = W|T|$ is the Polar Decomposition, there exist an orthonormal basis $\{e_n\}$ for $\ker |T|^\perp = \ker T^\perp$ and a sequence $\{\alpha_n\} \in l^1(\mathbb{C})$, s.t.

$$|T| = \sum_n \alpha_n e_n \otimes e_n$$

Then put $h_n = \sqrt{\alpha_n} e_n$ and $g_n = W h_n$. Since W is a partial isometry, $\{g_n\}$ is an orthonormal set. And clearly,

$$\|g_n\|^2 = \|h_n\|^2 = \alpha_n \Rightarrow \sum_n \|g_n\|^2 = \sum_n \|h_n\|^2 = \sum_n \alpha_n = \|T\|_1$$

and $T = \sum_n g_n \otimes h_n$. □

Remark. In fact, if g_n and h_n are chosen like above mention,

$$\|T\|_1 = \sum_n \|g_n\| \|h_n\|$$

Also, the structure of $\mathcal{B}_2(\mathcal{H})$ can be more explicit.

Theorem 14. 1) For $\mathcal{B}_2(\mathcal{H})$, define the inner product $\langle \cdot, \cdot \rangle$ as

$$\langle T, G \rangle = \text{tr } G^* T$$

Then $\|\cdot\|_2$ can be induced by this inner product, and moreover, $\mathcal{B}_2(\mathcal{H})$ is a Hilbert space with respect to this inner product.

2) If $A \in \mathcal{B}_0(\mathcal{H})$ and $\{\alpha_n\}$ are the eigenvalues of $|A|$, then $A \in \mathcal{B}_2(\mathcal{H})$ if and only if $(\alpha_n) \in l^2$. In this case, $\|A\|_1 = \sum \alpha_n^2$.

3) The Hilbert space $\mathcal{B}_2(\mathcal{H})$ containing $\mathcal{B}_{00}(\mathcal{H})$ as a dense subspace with respect to the norm $\|\cdot\|_2$.

5.3.2 Dual Spaces

In above subsection, we have seen $\mathcal{B}_1(\mathcal{H})$ is indeed a C^* -algebra, thus a Banach algebra with respect to the norm $\|\cdot\|_1$. Therefore, we can get the following important properties.

Theorem 1. For $A \in \mathcal{B}_1(\mathcal{H})$, define the functional

$$\begin{aligned} \Phi_A: \mathcal{B}_0(\mathcal{H}) &\longrightarrow \mathbb{C} \\ C &\longmapsto \text{tr } AC \end{aligned}$$

Then $\Phi_A \in \mathcal{B}_0(\mathcal{H})^*$. Therefore, there is a map

$$\begin{aligned} \rho: \mathcal{B}_1(\mathcal{H}) &\longrightarrow \mathcal{B}_0(\mathcal{H})^* \\ A &\longmapsto \Phi_A \end{aligned}$$

Moreover, ρ is an isomorphic isomorphism, i.e. $\mathcal{B}_1(\mathcal{H}) \cong \mathcal{B}_0(\mathcal{H})^*$.

Proof. Since $\mathcal{B}_0(\mathcal{H}) \subset \mathcal{B}_1(\mathcal{H})$, by above proposition, $\Phi_A \in \mathcal{B}_0(\mathcal{H})^*$. Moreover,

$$\sup \{ |\text{tr } AC| : \|C\|_1 \leq 1 \} \leq \|A\|_1 \Rightarrow \|\Phi_A\| \leq \|A\|_1$$

That means $\|\rho\| \leq 1$.

Check: ρ is surjective.

Let $\Phi \in \mathcal{B}_0(\mathcal{H})^*$ and define the sesquilinear map f on \mathcal{H} ,

$$f(g, h) = \Phi(g \otimes h), \text{ and } |f(g, h)| \leq \|\Phi\| \|g\| \|h\|$$

Therefore, there exists $A \in \mathcal{B}(\mathcal{H})$ s.t.

$$f(g, h) = \Phi(g \otimes h) = \langle Ag, h \rangle, \forall g, h \in \mathcal{H}$$

For any $C \in \mathcal{B}_{00}(\mathcal{H})$, then there are finite number $\{g_k\}$ and $\{h_k\}$ for $k = 1, \dots, n$ s.t. $C = \sum_{k=1}^n g_k \otimes h_k$

$$\begin{aligned} \Phi(C) &= \Phi\left(\sum_{k=1}^n g_k \otimes h_k\right) = \sum_{k=1}^n \langle Ag_k, h_k \rangle \\ &= \text{tr } AC = \Phi_A(C) \end{aligned}$$

Since $\mathcal{B}_{00}(\mathcal{H})$ is norm-dense in $\mathcal{B}_0(\mathcal{H})$ and Φ and Φ_A are bounded, $\Phi = \Phi_A$.

Check: ρ is isometric.

Let \mathcal{E} be an orthonormal basis and $A \in \mathcal{B}_1(\mathcal{H})$ with the Polar Decomposition $A = W|A|$, then for any finite subset E in \mathcal{E}

$$\begin{aligned} \|\Phi_A\| &\geq \left| \Phi\left(\sum_{e \in E} (e \otimes e) W^*\right) \right| = \left| \Phi_A\left(\sum_{e \in E} e \otimes We\right) \right| \\ &= \sum_{e \in E} |\langle Ae, We \rangle| \\ &= \sum_{e \in E} \langle |A| e, e \rangle \end{aligned}$$

Then $\|\Phi_A\| \geq \|A\|_1$ as $\{E\}$ can be an increasing net in \mathcal{E} . \square

Theorem 2. For $B \in \mathcal{B}(\mathcal{H})$, define the functional

$$\begin{aligned} \Psi_B: \quad \mathcal{B}_1(\mathcal{H}) &\longrightarrow \mathbb{C} \\ A &\longmapsto \operatorname{tr} AB \end{aligned}$$

Then $\Psi_B \in \mathcal{B}_1(\mathcal{H})^*$. Therefore, there is a map

$$\begin{aligned} \rho: \quad \mathcal{B}(\mathcal{H}) &\longrightarrow \mathcal{B}_1(\mathcal{H})^* \\ B &\longmapsto \Psi_B \end{aligned}$$

Moreover, ρ is an isomorphic isomorphism, i.e. $\mathcal{B}(\mathcal{H}) \cong \mathcal{B}_1(\mathcal{H})^*$.

Proof. Clearly, $\Psi_B \in \mathcal{B}_1(\mathcal{H})^*$ with $\|\Psi_B\| \leq \|B\|$.

Check: ρ is isometric.

For a given $\varepsilon > 0$, there exists a unit g s.t.

$$\|Bg\| > \|B\| - \varepsilon$$

And choosing a $h \in \mathcal{H}$, s.t. $\|Bg\| = \langle Bg, h \rangle$. Let $C = g \otimes h$, then $C \in \mathcal{B}_1(\mathcal{H})$ and $\|C\|_1 = 1$. Thus,

$$\|\Psi_B\| \geq |\operatorname{tr} BC| = \langle Bg, h \rangle = \|Bg\| > \|B\| - \varepsilon$$

Therefore, $\|\Psi_B\| = \|B\|$.

Check: ρ is surjective.

Let $\Psi \in \mathcal{B}_1(\mathcal{H})^*$. Similarly as above theorem, there is a $B \in \mathcal{B}(\mathcal{H})$ s.t. $\langle Bg, h \rangle = \Psi(g \otimes h)$ for all $g, h \in \mathcal{H}$. And also, $\Psi = \Psi_B$ on the $\mathcal{B}_{00}(\mathcal{H})$. By the extension, $\Psi = \Psi_B$. \square

Theorem 3. $(\mathcal{B}(\mathcal{H}), WOT)^* = (\mathcal{B}(\mathcal{H}), SOT)^* = \mathcal{B}_{00}(\mathcal{H})$.

Proof. By the **Proposition 2** in the subsection **3.3.3**, $L \in (\mathcal{B}(\mathcal{H}), WOT)^*$ can be expressed as

$$L(T) = \sum_{k=1}^n \langle Tg_k, h_k \rangle = \operatorname{tr} CT = \Psi_C(T), \text{ where } C = \sum_{k=1}^n g_k \otimes h_k$$

Thus, $(\mathcal{B}(\mathcal{H}), WOT)^* = (\mathcal{B}(\mathcal{H}), SOT)^* = \mathcal{B}_{00}(\mathcal{H})$. \square

5.3.3 Ultraweak Topology

By above subsection, we know that

$$\mathcal{B}(\mathcal{H}) \cong (\mathcal{B}_1(\mathcal{H}), \|\cdot\|_1)^*$$

Thus, $\mathcal{B}(\mathcal{H})$ can be equipped with the *weak** topology with respect to the $(\mathcal{B}_1(\mathcal{H}), \|\cdot\|_1)^*$. The subbasis of a point $A \in \mathcal{B}(\mathcal{H})$ is like,

$$V_\varepsilon(A) = \{ B \in \mathcal{B}(\mathcal{H}) : |\text{tr } T(A - B)| < \varepsilon, \forall T \in \mathcal{B}_1(\mathcal{H}) \}$$

for any $\varepsilon > 0$. This topology on $\mathcal{B}(\mathcal{H})$ is called the ultraweak topology or σ -*weak* topology. Therefore, for a net $\{A_\alpha\}$ in $\mathcal{B}(\mathcal{H})$, $A_\alpha \rightarrow A$ if and only of for any $T \in \mathcal{B}_1(\mathcal{H})$

$$\text{tr } TA_\alpha \rightarrow \text{tr } TA$$

Moreover, by the **Theorem 13** in above subsection, for any $T \in \mathcal{B}_1(\mathcal{H})$, there exist two sequences $\{g_n\}$ and $\{h_n\}$ s.t. $T = \sum_n g_n \otimes h_n$ and then

$$\text{tr } TA = \sum_n \langle Ag_n, h_n \rangle$$

Then, for a net $\{A_\alpha\}$ in $\mathcal{B}(\mathcal{H})$, $A_\alpha \rightarrow A$ if and only of for any two summable sequences $\{g_n\}$ and $\{h_n\}$ in \mathcal{H} ,

$$\sum_n \langle A_\alpha g_n, h_n \rangle \rightarrow \sum_n \langle Ag_n, h_n \rangle$$

Then there are some trivial properties of this topology.

Proposition 1. 1) If \mathcal{H} is separable, then the closed unit ball in $\mathcal{B}(\mathcal{H})$ with the ultraweak topology is a compact metric space.

2) The ultraweak topology and the WOT agree on bounded subsets of $\mathcal{B}(\mathcal{H})$.

3) A sequence in $\mathcal{B}(\mathcal{H})$ converges *weak** if and only if it converges WOT.

Proof. 1) is clearly by the form of subbasis of the ultraweak topology.

For 2), by above mention, the ultraweak topology is bigger than the WOT, that means the identity map

$$i : (\mathcal{B}(\mathcal{H}), \text{weak}^*) \longrightarrow (\mathcal{B}(\mathcal{H}), \text{WOT})$$

is continuous. When it restricts to any bounded subset \mathcal{S} , since \mathcal{S} is *weak**-compact by 1), i is a homeomorphism.

By the Principle of Uniform Boundedness, WOT -convergent sequence is bounded, thus by 2), it is $weak^*$ -convergent. And the converse is trivial. \square

Now, we want to find more relations between the ultraweak topology and the WOT . Firstly, the inflation can provide us some information about that.

Proposition 2. *If $1 \leq d \leq \infty$ and $C \in \mathcal{B}_1(\mathcal{H}^{(d)})$ with the matrix representation $C = [C_{jk}]$, then*

1) $K = \sum_k C_{kk}$ converges $weak^*$ in $\mathcal{B}_1(\mathcal{H})$ and

$$\left\| \sum_k C_{kk} \right\| \leq \|C\|_1$$

2) if $T \in \mathcal{B}(\mathcal{H})$ and K is as 1), then

$$\text{tr } T^{(d)} C = \text{tr } TK = \sum_{k=1}^d \text{tr } TC_{kk}$$

and this converges absolutely for $d = \infty$.

Proof. We just need to prove this proposition in the infinite case.

For 1), let P_k be the projection from $\mathcal{H}^{(\infty)}$ onto the k -th coordinate space. Then $C_{kk} = P_k C P_k$. And clearly, for $1 \leq n < m < \infty$, $\sum_{k=n}^m C_{kk} \in \mathcal{B}_1(\mathcal{H})$, if $L \in \mathcal{B}_0(\mathcal{H}^{(\infty)})$,

$$\begin{aligned} \left| \text{tr} \left(L \sum_{k=n}^m C_{kk} \right) \right| &= \left| \sum_{k=n}^m \text{tr} (L P_k C P_k) \right| = \left| \sum_{k=n}^m \text{tr} (C P_k L P_k) \right| \\ &= \left| \text{tr} \left(C \sum_{k=n}^m P_k L P_k \right) \right| \\ &\leq \|C\|_1 \left\| \sum_{k=n}^m P_k L P_k \right\| \\ &= \|C\|_1 \sup \{ \|P_k L P_k\| : n \leq k \leq m \} \\ &\leq \|C\|_1 \|L\| \end{aligned}$$

Then by the **Proposition 3** in the subsection **3.2.1**, (Confused!)

$$\left\| \sum_{k=n}^m C_{kk} \right\|_1 = \sup \left\{ \left| \text{tr} \left(L \sum_{k=n}^m C_{kk} \right) \right| : \|L\| \leq 1, L \in \mathcal{B}_0(\mathcal{H}^{(\infty)}) \right\}$$

And since $L \in \mathcal{B}_0(\mathcal{H}^{(\infty)})$, $\|P_k L P_k\| \rightarrow 0$. Therefore, 1) holds.

For 2), let $\{e_i: i \in I\}$ be a orthonormal basis of \mathcal{H} and $\{e_{ki}: 1 \leq k \leq \infty, i \in I\}$ be a orthonormal basis of $\mathcal{H}^{(\infty)}$, then for any $T \in \mathcal{B}(\mathcal{H})$, $TC \in \mathcal{B}_1(\mathcal{H}^{(\infty)})$ and

$$\text{tr}(T^{(\infty)}C) = \sum_{i \in I} \sum_{k=1}^{\infty} \langle T^{(\infty)}C e_{ki}, e_{ki} \rangle$$

Fixing k , we have that

$$\langle T^{(\infty)}C e_{ki}, e_{ki} \rangle = \langle TC_{kk}e_i, e_i \rangle = \text{tr}((e_i \otimes e_i)TC_{kk})$$

By the *weak**-convergence of $K = \sum_k C_{kk}$,

$$\begin{aligned} \langle TK e_i, e_i \rangle &= \sum_{k=1}^{\infty} \text{tr}((e_i \otimes e_i)TC_{kk}) \\ &= \sum_{k=1}^{\infty} \langle TC_{kk}e_i, e_i \rangle \\ &= \sum_{k=1}^{\infty} \langle T^{(\infty)}C e_{ki}, e_{ki} \rangle \end{aligned}$$

Since above series is absolutely convergent,

$$\text{tr}(T^{(\infty)}C) = \sum_{k=1}^{\infty} \sum_{i \in I} \langle TC_{kk}e_i, e_i \rangle = \sum_{k=1}^{\infty} \text{tr}(TC_{kk}) = \text{tr}TK$$

□

Theorem 3. 1) If $C \in \mathcal{B}_1(\mathcal{H})$ with rank at most d , where $1 \leq d \leq \infty$, then there are $g, h \in \mathcal{H}^{(d)}$ s.t. $\|g\|^2 = \|h\|^2 \leq \|C\|_1$ and

$$\langle T^{(d)}g, h \rangle = \text{tr}TC, \quad \forall T \in \mathcal{B}(\mathcal{H})$$

2) If $1 \leq d \leq \infty$ and $g, h \in \mathcal{H}^{(d)}$, then there is a $C \in \mathcal{B}_1(\mathcal{H})$ with rank at most d s.t.

$$\langle T^{(d)}g, h \rangle = \text{tr}TC, \quad \forall T \in \mathcal{B}(\mathcal{H})$$

3) If $C \in \mathcal{B}_1(\mathcal{H}^{(\infty)})$, then there are $g, h \in \mathcal{H}^{(\infty)}$ s.t. $\|g\|^2 = \|h\|^2 \leq \|C\|_1$ and

$$\langle T^{(\infty)}g, h \rangle = \text{tr}T^{(\infty)}C, \quad \forall T \in \mathcal{B}(\mathcal{H})$$

4) If $g, h \in \mathcal{H}^{(\infty)}$, then there is a $C \in \mathcal{B}_1(\mathcal{H}^{(\infty)})$ s.t. $\|g\|^2 = \|h\|^2 \leq \|C\|_1$ and

$$\langle T^{(\infty)}g, h \rangle = \text{tr}T^{(\infty)}C, \quad \forall T \in \mathcal{B}(\mathcal{H})$$

Proof. Just need to prove the case when $d = \infty$.

For 1), let $C = \sum_k g_k \otimes h_k$ with $\|g_k\|^2 = \|h_k\|^2 = \|C\|_1$. Then $g = (g_k)$ and $h = (h_k)$ are in $\mathcal{H}^{(\infty)}$, and it is clear that

$$\langle T^{(\infty)}g, h \rangle = \text{tr } TC, \quad \forall T \in \mathcal{B}(\mathcal{H})$$

And the converse is similar.

If $C = (C_{jk}) \in \mathcal{B}_1(\mathcal{H}^{(\infty)})$, then by above proposition, $K = \sum_k C_{kk} \in \mathcal{B}_1(\mathcal{H})$ and $\|K\|_1 \leq \|C\|_1$. Moreover, $\text{tr}(T^{(\infty)}C) = \text{tr}(TK)$. Then by 1), there are $g = (g_k)$ and $h = (h_k)$ in $\mathcal{H}^{(\infty)}$ with $\|g_k\|^2 = \|h_k\|^2 = \|K\|_1 \leq \|C\|_1$ s.t.

$$\text{tr}(T^{(\infty)}C) = \text{tr}(TK) = \langle T^{(\infty)}g, h \rangle.$$

And 4) is also trivial. \square

Remark. In fact, it is the generalization of the result combining the **Proposition 11** and the **Theorem 13** in the subsection **5.3.1**.

Then, we can see one more relation between the ultraweak topology and the *WOT*.

Theorem 4. *If $\{T_i\}$ is a net in $\mathcal{B}(\mathcal{H})$, then the following statements are equivalent.*

- 1) $T_i \rightarrow 0$ in weak*.
- 2) $T_i^\infty \rightarrow 0$ in weak*.
- 3) $T_i^\infty \rightarrow 0$ in *WOT*.

Proof. If $C \in \mathcal{B}_1(\mathcal{H}^{(\infty)})$ and $K = \sum_k C_{kk}$, then by above theorem,

$$\text{tr } T_i^{(\infty)} C = \text{tr}(T_i K)$$

Therefore, 1) implies 2). Since the the ultraweak topology and the *WOT* agree on any bounded subset, 2) implies 3). Let $K \in \mathcal{B}_1(\mathcal{H})$, then there are $h, g \in \mathcal{H}^\infty$ s.t. $\text{tr}(KT_i) = \langle T_i^\infty f, g \rangle$. Thus, 3) implies 1). \square

Corollary 5. 1) *If $\mathcal{S} \subset \mathcal{B}(\mathcal{H})$ is a subset and $\mathcal{S}_1 = \overline{\mathcal{S}}^{wk*}$, then*

$$\overline{\{S^{(\infty)} : S \in \mathcal{S}\}}^{WOT} = \overline{\{S^{(\infty)} : S \in \mathcal{S}\}}^{wk*} = \mathcal{S}_1^{(\infty)}$$

Therefore, \mathcal{S} is weak-closed if and only if $\mathcal{S}^{(\infty)}$ is *WOT*-closed.*

2) *If \mathfrak{A} is a C^* -subalgebra of $\mathcal{B}(\mathcal{H})$, then*

$$\overline{\mathfrak{A}}^{wk*} = \mathfrak{A}''$$

That means any von Neumann algebra is weak-closed.*

Proof. 1) is a direct result from above theorem.

For 2), let $\mathfrak{A}_1 = \overline{\mathfrak{A}}^{wk*}$. Then by 1) and the Double Commutant Theorem,

$$\mathfrak{A}_1^{(\infty)} = \{\mathfrak{A}^{(\infty)}\}'' = \{\mathfrak{A}''\}^{(\infty)}$$

Thus, $\mathfrak{A}_1 = \mathfrak{A}''$. □

Remark. Conversely, an algebra \mathfrak{A} is *weak**-closed, it is trivially *WOT*-closed. Therefore, we have another powerful topology to research von Neumann algebras and it is very important.

5.4 Fredholm Theory

Let $\pi: \mathcal{B}(\mathcal{H}) \rightarrow \mathcal{B}(\mathcal{H})/\mathcal{B}_0(\mathcal{H})$ be the quotient map. Then

$$\pi(A) = A + \mathcal{B}_0(\mathcal{H})$$

And the C^* -algebra $\mathcal{B}(\mathcal{H})/\mathcal{B}_0(\mathcal{H})$ is called the Calkin algebra. Therefore, in order to research the compact perturbations, the Calkin algebra plays an important role.

5.4.1 Fredholm Operators

Definition 1. Let $A \in \mathcal{B}(\mathcal{H})$ and $\pi: \mathcal{B}(\mathcal{H}) \rightarrow \mathcal{B}(\mathcal{H})/\mathcal{B}_0(\mathcal{H})$ be the quotient map. If $\pi(A)$ is left-invertible (or right) in the Calkin algebra, then A is called a left (or right) semi-Fredholm operator. If $\pi(A)$ is invertible in the Calkin algebra, then A is called a Fredholm operator.

There are some equivalent conditions of the Fredholm operators.

Theorem 2. If $A \in \mathcal{B}(\mathcal{H})$, then the following statements are equivalent.

- 1) A is left semi-Fredholm.
- 2) $\text{ran } A$ is closed and $\dim \ker A < \infty$.
- 3) There is a $B \in \mathcal{B}(\mathcal{H})$ and a finite rank $F \in \mathcal{B}(\mathcal{H})$ s.t.

$$BA = 1 + F$$

Proof. 1) \Rightarrow 2): There is a B s.t. $\pi(B)\pi(A) = 1$, i.e. $\pi(BA - 1) = 0$. That means there is a $K \in \mathcal{B}_0(\mathcal{H})$ s.t. $BA = 1 + K$. Then, by the **Theorem 8** in the subsection **5.1.2**,

$$\ker A \subset \ker BA = \ker (1 + K) < \infty$$

And $\text{ran } BA = \text{ran } (1 + K)$ is closed. By the **Lemma 3** in the subsection **2.4.4**, there is a constant $C > 0$ s.t. for any $h \in (\ker BA)^\perp$,

$$C \|h\| \leq \|BAh\| \leq \|B\| \|Ah\| \Rightarrow \|Ah\| > C' \|h\|$$

Therefore, $A((\ker BA)^\perp)$ is closed. Because

$$\text{ran } A = A((\ker BA)^\perp) + A(\ker BA)$$

and $\ker BA < \infty$, then $\text{ran } A$ is closed.

2) \Rightarrow 3): Define $A_1: (\ker A)^\perp \rightarrow \text{ran } A$, then A_1 is invertible by the Inverse Mapping Theorem. Let F be the projection onto $\ker A$ and P be the projection onto $\text{ran } A$, then

$$BA = 1 - F, \text{ where } B = A_1^{-1}P$$

3) implies 1) trivially. □

Above theorem provides the information of the range and the kernel of Fredholm operators. Besides that, we can also find the equivalent conditions of Fredholm operators from another point of view.

Theorem 3. *If $A \in \mathcal{B}(\mathcal{H})$, then the following statements are equivalent.*

- 1) A is left semi-Fredholm.
- 2) There is no unit sequence $\{h_n\}$ s.t. $h_n \rightarrow 0$ weakly and $Ah_n \rightarrow 0$ in norm.
- 3) There is no orthonormal sequence $\{e_n\}$ s.t. $Ae_n \rightarrow 0$ in norm.
- 4) There is a $\delta > 0$ s.t. $\{h \in \mathcal{H}: \|Ah\| \leq \delta \|h\|\}$ contains no infinite dimensional manifold.
- 5) Let $|A| = \int t dE(t)$, then there is a $\delta > 0$ s.t. $E([0, \delta])\mathcal{H}$ is finite dimensional.

Proof. 1) \Rightarrow 2): Let $B \in \mathcal{B}(\mathcal{H})$ and $T \in \mathcal{B}_0(\mathcal{H})$ s.t. $BA = 1 + K$. For a unit sequence $\{h_n\}$ with $h_n \rightarrow 0$ weakly, since K is completely continuous,

$$|1 - \|BAh_n\|| = |||h_n\| - \|BAh_n\|| \leq \|Kh_n\| \rightarrow 0$$

Therefore, $\|BAh_n\| \rightarrow 1$, i.e. Ah_n cannot converge to 0 in norm.

2) \Rightarrow 3): Orthonormal sequences always converge weakly 0.

3) \Rightarrow 4): Assume it is not true. For any n , there is a manifold \mathcal{M}_n s.t.

$$\|Ah\| \leq \frac{1}{n} \|h\|, \quad \forall h \in \mathcal{M}_n$$

Let $e_k \in \mathcal{M}_k$ for $1 \leq k \leq n$ and $\{e_k\}_{k=1}^n$ be orthonormal and E be the projection onto $\text{span}\{e_k\}_{k=1}^n$. Since \mathcal{M}_{n+1} is infinite dimensional, $\mathcal{M}_{n+1} \cap E^\perp \neq \{0\}$. Thus, there is a unit $e_{n+1} \in \mathcal{M}_{n+1}$ s.t. $e_{n+1} \perp E$. But $Ae_n \rightarrow 0$ in norm contradicted to 3).

4) \Rightarrow 5): If $\delta > 0$ and $h \in E([0, \delta])\mathcal{H}$, then

$$\begin{aligned} \|Ah\|^2 &= \langle A^*Ah, h \rangle = \langle |A|^2 h, h \rangle \\ &= \int_0^\delta t^2 dE_{hh}(t) \\ &\leq \delta^2 E_{hh}([0, \delta]) = \delta^2 \|h\|^2 \end{aligned}$$

That means

$$E([0, \delta])\mathcal{H} \subset \{h \in \mathcal{H} : \|Ah\| \leq \delta \|h\|\}$$

5) \Rightarrow 1): Let $\mathcal{M} = (E([0, \delta])\mathcal{H})^\perp$. Then $|A|$ maps \mathcal{M} onto \mathcal{M} bijectively. Put

$$B_1 = \int_\delta^\infty t^{-1} dE(t)$$

then $B_1 |A| = 1|_{\mathcal{M}}$. Let $A = U|A|$. Since $\mathcal{M} \subset \text{ran } |A|$, U maps \mathcal{M} isometrically to $U(\mathcal{M})$. Then there is a V s.t. $VU = 1|_{\mathcal{M}}$ and $V|_{U(\mathcal{M})^\perp} = 0$. Put $B = B_1 V$, then for any $h \in \mathcal{M}$,

$$BAh = B_1 VU |A| h = h$$

If $h \in \mathcal{M}^\perp$, we know that $BAh = 0$. Thus,

$$BA = E((\delta, \infty)) = 1 - E([0, \delta])$$

Thus, by above theorem, A is left semi-Fredholm. \square

Theorem 4. $A \in \mathcal{B}(\mathcal{H})$ is left semi-Fredholm operator if and only if for any $K \in \mathcal{B}_0(\mathcal{H})$,

$$\dim \ker (A + K) < \infty$$

Proof. If there is $B \in \mathcal{B}(\mathcal{H})$ and $L \in \mathcal{B}_0(\mathcal{H})$ s.t. $BA = 1 + L$. If $K \in \mathcal{B}_0(\mathcal{H})$,

$$B(A + K) = 1 + (L + BK)$$

Thus $A + K$ is left semi-Fredholm, then $\dim \ker (A + K) < \infty$.

Conversely, if there is a orthonormal sequence $\{e_n\}$ s.t. $Ae_n \rightarrow 0$ in norm.

By passing a subsequence, it may assume that $\sum \|Ae_n\|^2 < \infty$. Thus for any $h \in \mathcal{H}$,

$$\begin{aligned} \sum |\langle h, e_n \rangle| \|Ae_n\| &\leq \left(\sum |\langle h, e_n \rangle|^2 \right)^{\frac{1}{2}} \left(\sum \|Ae_n\|^2 \right)^{\frac{1}{2}} \\ &\leq C \|h\| \end{aligned}$$

Thus $Kh = \sum \langle h, e_n \rangle Ae_n$ is a bounded operator and clearly it can be approximated by finite rank operators $K_n h = \sum_{k=1}^n \langle h, e_k \rangle Ae_k$, thus $K \in \mathcal{B}_0(\mathcal{H})$. But $(A - K)e_n = 0$ for all n , that means $\dim \ker(A - K) = \infty$. \square

Corollary 5. $A \in \mathcal{B}(\mathcal{H})$ is Fredholm if and only if $\text{ran } A$ is closed and both $\ker A$ and $\ker A^*$ are finite dimensional.

5.4.2 Fredholm Index

Definition 1. If A is a semi-Fredholm operator, then the Fredholm index of A is defined as

$$\text{ind } A = \dim \ker A - \dim \ker A^*$$

Remark. In fact, $\text{ind } A \in \mathbb{Z} \cup \{\pm\infty\}$. If A is Fredholm, $\text{ind } A \in \mathbb{Z}$. And clearly, $\text{ind } (\lambda + K) = 0$ for all nonzero $\lambda \in \mathbb{C}$ and $K \in \mathcal{B}_0(\mathcal{H})$.

Proposition 2. 1) If A is a semi-Fredholm, then so is A^* and

$$\text{ind } A^* = -\text{ind } A$$

2) If N is normal, then N is semi-Fredholm if and only if N is Fredholm, in which case $\text{ind } N = 0$.

3) If A and B are Fredholm, then $A \oplus B$ is Fredholm and

$$\text{ind } A \oplus B = \text{ind } A + \text{ind } B$$

Proof. 1) and 3) is trivial and 2) is since $\|N\| = \|N^*\|$. \square

In fact, for $A \in \mathcal{B}(\mathcal{H}, \mathcal{H}')$, we can similarly define the Fredholm operator and the Fredholm index, and all of above mention are valid.

Proposition 3. If \mathcal{H} and \mathcal{H}' are finite dimensional Hilbert spaces and $A \in \mathcal{B}(\mathcal{H}, \mathcal{H}')$, then A is Fredholm and $\text{ind } A = \dim \mathcal{H} - \dim \mathcal{H}'$.

Proof. The fact that A is Fredholm is trivial by the definition.

$$\begin{aligned} \text{ind } A &= \dim \ker A - \dim (\text{ran } A)^\perp \\ &= \dim \mathcal{H} - \dim \text{ran } A - (\dim \mathcal{H}' - \dim \text{ran } A) \\ &= \dim \mathcal{H} - \dim \mathcal{H}' \end{aligned}$$

\square

Lemma 4. Let $A \in \mathcal{B}(\mathcal{H}, \mathcal{H}')$ and $\mathcal{H} = \mathcal{M} \oplus \mathcal{N}$ and $\mathcal{H}' = \mathcal{M}' \oplus \mathcal{N}'$. Suppose A has the form

$$A = \begin{pmatrix} A_1 & X \\ 0 & A_2 \end{pmatrix}$$

If $\dim \mathcal{N} < \infty$ and $\dim \mathcal{N}' < \infty$ and A_1 is invertible, then A is Fredholm and $\text{ind } A = \dim \mathcal{N} - \dim \mathcal{N}'$.

Proof. Just need to show $\ker A^* = \ker A_2^*$ and $\dim \ker A = \dim \ker A_2$, then since $A_2 \in \mathcal{B}(\mathcal{N}) = \mathcal{B}_0(\mathcal{N})$,

$$\text{ind } A = \dim \ker A - \dim \ker A^* = \text{ind } A_2 = \dim \mathcal{N} - \dim \mathcal{N}'$$

But $\ker A^* = \ker A_2^*$ and $\dim \ker A = \dim \ker A_2$ are clearly because A_1 is invertible. \square

Theorem 5. If $A, B \in \mathcal{B}(\mathcal{H})$ are Fredholm operator, then BA is also a Fredholm operator and

$$\text{ind } BA = \text{ind } A + \text{ind } B$$

Proof. There are $X, Y \in \mathcal{B}(\mathcal{H})$ s.t. $XA = 1 + K$ and $YB = 1 + K'$ for some $K, K' \in \mathcal{B}_0(\mathcal{H})$. Then

$$(XY)(BA) = X(1 + K')A = 1 + (K + XK'A)$$

Therefore, BA is left semi-Fredholm. Similarly, BA is right semi-Fredholm. Then for the index, let

$$\begin{aligned} \mathcal{M}' &= \text{ran } A \cap (\ker B)^\perp, & \mathcal{N}' &= \mathcal{M}'^\perp \\ \mathcal{M} &= A^{-1}(\mathcal{M}') \cap (\ker A)^\perp, & \mathcal{N} &= \mathcal{M}^\perp \\ \mathcal{M}'' &= B(\mathcal{M}'), & \mathcal{N}'' &= \mathcal{M}''^\perp \end{aligned}$$

Therefore,

$$\mathcal{H} = \mathcal{M} \oplus \mathcal{N} = \mathcal{M}' \oplus \mathcal{N}' = \mathcal{M}'' \oplus \mathcal{N}''$$

And since $A(\mathcal{M}) = \mathcal{M}'$, we can see

$$\begin{aligned} A &= \begin{pmatrix} A_1 & X \\ 0 & A_2 \end{pmatrix} : \mathcal{M} \oplus \mathcal{N} \longrightarrow \mathcal{M}' \oplus \mathcal{N}' \\ B &= \begin{pmatrix} B_1 & Y \\ 0 & B_2 \end{pmatrix} : \mathcal{M}' \oplus \mathcal{N}' \longrightarrow \mathcal{M}'' \oplus \mathcal{N}'' \\ BA &= \begin{pmatrix} B_1 A_1 & Z \\ 0 & B_2 A_2 \end{pmatrix} : \mathcal{M} \oplus \mathcal{N} \longrightarrow \mathcal{M}'' \oplus \mathcal{N}'' \end{aligned}$$

It can be seen that \mathcal{N} , \mathcal{N}' and \mathcal{N}'' are finite dimensional. Moreover, A_1 , B_1 and $B_1 A_1$ are invertible. Then using above lemma,

$$\begin{aligned}\operatorname{ind} BA &= \dim \mathcal{N} - \dim \mathcal{N}'' \\ &= \dim \mathcal{N} - \dim \mathcal{N}' + \dim \mathcal{N}' - \dim \mathcal{N}'' \\ &= \operatorname{ind} B + \operatorname{ind} A\end{aligned}$$

□

Remark. In fact, above theorem can also be valid if A and B are semi-Fredholm.

Theorem 6. *If $A \in \mathcal{B}(\mathcal{H})$ is a Fredholm operator, then for any $K \in \mathcal{B}(\mathcal{H})$, $A + K$ is Fredholm and $\operatorname{ind}(A + K) = \operatorname{ind} A$.*

Proof. $A + K$ is Fredholm is trivial. Let X be Fredholm and $L \in \mathcal{B}_0(\mathcal{H})$, s.t. $XA = 1 + L$, then

$$0 = \operatorname{ind}(1 + L) = \operatorname{ind} XA = \operatorname{ind} X + \operatorname{ind} A$$

And since $X(A + K) = 1 + (L + XK)$,

$$\operatorname{ind} X(A + K) = \operatorname{ind}(A + K) + \operatorname{ind} X = 0$$

Therefore, $\operatorname{ind}(A + K) = \operatorname{ind} A$. □

Now, we can see that the index of a Fredholm is invariant under a small perturbation.

Theorem 7. *If $A \in \mathcal{B}(\mathcal{H})$ is a Fredholm operator, then there is an $\varepsilon > 0$ s.t. for any $Y \in \mathcal{B}(\mathcal{H})$ with $\|Y\| < \varepsilon$, then $A + Y$ is Fredholm and $\operatorname{ind} A + Y = \operatorname{ind} A$.*

Proof. Let $\mathcal{H} = \ker A^\perp \oplus \ker A = \operatorname{ran} A \oplus \ker A^*$. Then,

$$A = \begin{pmatrix} A_1 & 0 \\ 0 & 0 \end{pmatrix}$$

And A_1 is invertible since $\operatorname{ran} A$ is closed. By the perturbation of inverse, there is a $\varepsilon > 0$ s.t. if $\|Y_1\| < \varepsilon$, $A_1 + Y_1$ is invertible. If $Y \in \mathcal{B}(\mathcal{H})$ with $\|Y\| < \varepsilon$, and

$$Y = \begin{pmatrix} Y_1 & Y_2 \\ Y_3 & Y_4 \end{pmatrix}$$

Thus $A_1 + Y_1$ is invertible and

$$A + Y = \begin{pmatrix} A_1 + Y_1 & Y_2 \\ Y_3 & Y_4 \end{pmatrix} = \begin{pmatrix} A_1 + Y_1 & 0 \\ 0 & 0 \end{pmatrix} + \begin{pmatrix} 0 & Y_2 \\ Y_3 & Y_4 \end{pmatrix}$$

The first part is Fredholm and the second part is finite rank. Thus $A + Y$ is Fredholm, and $\operatorname{ind} A + Y = \operatorname{ind} A$ by above theorem. □

Corollary 8. Let \mathcal{F} denote the set of all Fredholm operators. Then the map

$$\text{ind} : (\mathcal{F}, \|\cdot\|) \longrightarrow \mathbb{Z}$$

is continuous with respect to the discrete topology on \mathbb{Z} .

5.4.3 Essential Spectrum

Let $\pi : \mathcal{B}(\mathcal{H}) \rightarrow \mathcal{B}(\mathcal{H})/\mathcal{B}_0(\mathcal{H})$ be the quotient map. We have know that the inverse of the invertible element in the Calkin algebra is Fredholm. Thus, the spectrum of the element in the Calkin algebra may be interesting.

Definition 1. If $A \in \mathcal{B}(\mathcal{H})$, the essential spectrum of A is defined as

$$\sigma_e(A) = \sigma(\pi(A))$$

Similarly, $\sigma_{le} = \sigma_l(\pi(A))$ and $\sigma_{re} = \sigma_r(\pi(A))$.

By the properties of the general spectrum and the Fredholm operators, we have following properties of the essential spectrum.

Proposition 2. Let $A \in \mathcal{B}(\mathcal{H})$.

- 1) $\sigma_{le}(A) \subset \sigma_l(A)$ and $\sigma_{re}(A) \subset \sigma_r(A)$ and $\sigma_e(A) \subset \sigma(A)$.
- 2) $\sigma_{le}(A)$, $\sigma_{re}(A)$ and $\sigma_e(A)$ are compact.
- 3) For any $K \in \mathcal{B}_0(\mathcal{H})$, $A+K$ and A have same left (or right) essential spectrum.
- 4) $\lambda \in \sigma_{le}(A)$ is and only if $\dim \ker(A - \lambda) = \infty$ or $\text{ran}(A - \lambda)$ is not closed.
- 5) $\lambda \in \sigma_{re}(A)$ is and only if $\dim \text{ran}(A - \lambda)^\perp = \infty$ or $\text{ran}(A - \lambda)$ is not closed.

Proposition 3. If $A \in \mathcal{B}(\mathcal{H})$, then

$$\sigma_{ap}(A) = \sigma_{le}(A) \cup \{\lambda \in \sigma_p(A) : \dim \ker(A - \lambda) < \infty\}$$

Proof. If $\lambda \in \sigma_{ap}(A)$, then by the **Proposition 5** and **8** in the subsection **3.3.1**, we know that either $\text{ran}(A - \lambda)$ is not closed or $\ker(A - \lambda) \neq \{0\}$. If $\text{ran}(A - \lambda)$ is not closed or $\dim \ker(A - \lambda) = \infty$, $\lambda \in \sigma_{le}(A)$. The converse is similar. \square

For normal operators, by above proposition, there are more interesting results about the essential spectrum.

Lemma 4. If N is a normal operator and $\lambda \in \sigma(N)$, then $\text{ran}(N - \lambda)$ is closed if and only if λ is not a limit point of $\sigma(N)$.

Proof. Assume λ is isolated in $\sigma(N)$. Let $X = \sigma(N) \setminus \{\lambda\}$ and $\mathcal{H}_1 = E(X)\mathcal{H}$. Hence, $(N - \lambda)\mathcal{H}_1$ is closed as $\sigma(N|_{\mathcal{H}_1}) = X$. Since

$$\mathcal{H}_1^\perp = E(\{\lambda\})\mathcal{H} = \ker(A - \lambda)$$

$\text{ran}(N - \lambda) = (N - \lambda)\mathcal{H}_1$. Therefore, $\text{ran}(N - \lambda)$ is closed.

Conversely, assume λ is not isolated in $\sigma(N)$. Then there is a positive sequence $\{r_n\}$ s.t. $r_n \rightarrow 0$ decreasingly. Each open set

$$A_n = \{z \in \mathbb{C} : r_{n+1} < |z - \lambda| < r_n\}$$

has non-empty intersection with $\sigma(N)$, i.e. $E(A_n)\mathcal{H} \neq \{0\}$ for all n . Let $e_n \in A_n$ be the unit element and $e_n \in \ker(N - \lambda)^\perp$ and

$$\|(N - \lambda)e_n\|^2 = \int_{A_n} |z - \lambda|^2 dE_{e_n, e_n}(z) \leq r_n^2 \rightarrow 0$$

Therefore,

$$\inf \{\|(N - \lambda)h\| : \|h\| = 1, h \in \ker(N - \lambda)^\perp\} = 0$$

i.e. $\text{ran}(N - \lambda)$ is not closed. \square

Theorem 5. *If N is normal, then $\sigma_e(N) = \sigma_{le}(N) = \sigma_{re}(N)$ and*

$$\sigma(N) \setminus \sigma_e(N) = \{\lambda \in \sigma(N) : \lambda \text{ is isolated and } \dim \ker(N - \lambda) < \infty\}$$

Proof. The first part can be obtained by applying the **Corollary 6** in the subsection **3.3.1** to the Calkin algebra.

If λ is isolated in $\sigma(N)$, then $\text{ran}(N - \lambda)$ is closed. Thus,

$$\dim \ker(N - \lambda) < \infty \Rightarrow \lambda \notin \sigma_{le}(N)$$

The converse is similar. \square

Let $\gamma(A) = \inf \{\|Ah\| : \|h\| = 1, h \perp \ker A\} = \chi(0)$. In fact, $\gamma(A) = \gamma(A^*)$. Moreover,

$$\gamma(A) \text{dist}(h, \ker A) \leq \|Ah\|$$

Lemma 6. *If \mathcal{M}, \mathcal{N} are closed subspace of \mathcal{H} and \mathcal{N} is finite dimensional and $\dim \mathcal{M} > \dim \mathcal{N}$, then there is a non-zero vector $m \in \mathcal{M}$ s.t.*

$$\|m\| = \text{dist}(m, \mathcal{N})$$

Proof. Let P be the projection onto \mathcal{M} , so $\dim P(\mathcal{N}) \leq \dim \mathcal{N} < \dim \mathcal{M}$. Thus $P(\mathcal{N})$ is a proper subspace of \mathcal{M} . Let $m \in \mathcal{M} \cap P(\mathcal{N})^\perp$. For any $n \in \mathcal{N}$,

$$0 = \langle Pn, m \rangle = \langle n, Pm \rangle = \langle n, m \rangle$$

Thus $m \perp \mathcal{N}$, i.e. $\|m\| = \text{dist}(m, \mathcal{N})$. \square

Lemma 7. *If A is left semi-Fredholm and $B \in \mathcal{B}(\mathcal{H})$ s.t. $\|B\| < \gamma(A)$, then $A + B$ is left semi-Fredholm and*

- 1) $\dim \ker (A + B) \leq \dim \ker A$.
- 2) $\dim \operatorname{ran} (A + B)^\perp \leq \dim \operatorname{ran} A^\perp$.

Proof. Check: $\{h: \|(A + B)h\| \leq \delta \|h\|\}$ for some $\delta < \gamma - \|B\|$ contains no infinite dimensional subspace.

If there is an infinite dimensional subspace, then it would contain a finite dimensional subspace \mathcal{M} with $\dim \mathcal{M} > \dim \ker A$. Therefore, there is a $h \in \ker A$ s.t. $\|h\| = \operatorname{dist}(h, \ker A)$. Then

$$\begin{aligned} \gamma(A) \|h\| &= \gamma(A) \operatorname{dist}(h, \ker A) \leq \|Ah\| \\ &\leq \|(A + B)h\| + \|Bh\| \leq (\delta + \|B\|) \|h\| \\ &< \gamma(A) \|h\| \end{aligned}$$

which is a contradiction. Therefore, $A + B$ is a left semi-Fredholm.

Check: $\dim \ker (A + B) \leq \dim \ker A$.

If $h \in \ker (A + B)$, then $Ah = -Bh$. Then

$$\gamma(A) \operatorname{dist}(h, \ker A) \leq \|Bh\| \leq \|B\| \|h\| < \gamma(A) \|h\|$$

Therefore, $\operatorname{dist}(h, \ker A) < \|h\|$. By above lemma,

$$\dim \ker (A + B) \leq \dim \ker A$$

Finally, if $\dim \operatorname{ran} A^\perp = \infty$, 2) holds clearly. But if $\dim \operatorname{ran} A^\perp < \infty$, A is also right semi-Fredholm, i.e. A^* is left semi-Fredholm, and by 1), 2) holds. \square

Remark. Clearly, these results can also be valid for right semi-Fredholm and Fredholm.

Proposition 8. *If A is semi-Fredholm and either $\ker A = \{0\}$ or $\operatorname{ran} A = \mathcal{H}$, then there is a $\delta > 0$ s.t. if $\|B - A\| < \delta$, the*

$$\dim \ker A = \dim \ker B, \quad \dim \operatorname{ran} A = \dim \operatorname{ran} B$$

Proof. In fact, there is a $\delta > 0$ s.t. if $\|B - A\| < \delta$, then $\operatorname{ind} A = \operatorname{ind} B$ and

$$\dim \ker B \leq \dim \ker A, \quad \dim \operatorname{ran} B^\perp \leq \dim \operatorname{ran} A^\perp$$

Since one of these is 0, the results are true. \square

Theorem 9. *If $\lambda \notin \sigma_{le}(A) \cap \sigma_{re}(A)$, then there is a $\delta > 0$ s.t. $\dim \ker (A - \mu)$ and $\dim \operatorname{ran} (A - \mu)^\perp$ are constant for $0 < |\lambda - \mu| < \delta$.*

Proof. Assume $\lambda = 0$ and $\ker A$ is finite dimensional. And A^n is also left semi-Fredholm for all n . Let $\mathcal{M}_n = \text{ran } A^n$. Note $\mathcal{M}_{n+1} \subset \mathcal{M}_n$ and $\mathfrak{A}\mathcal{M}_n = \mathcal{M}_{n+1}$. Let $\mathcal{M} = \bigcap \mathcal{M}_n$ and $B = A|_{\mathcal{M}}$.

Claim: $B(\mathcal{M}) = \mathcal{M}$.

Since $\ker A$ is finite dimensional and $\{\mathcal{M}_n\}$ is decreasing, there is a m s.t. for all $n \geq m$, $\mathcal{M}_n \cap \ker A = \mathcal{M}_m \cap \ker A$. For any $h \in \mathcal{M}$ and $n \geq m$, there is $f_n \in \mathcal{M}_n$ s.t. $h = Af_n$.

$$f_m - f_n \in \mathcal{M}_n \cap \ker A = \mathcal{M}_m \cap \ker A$$

Therefore, $f_m \in \mathcal{M}_n$ for all $n \geq m$. Thus $f_m \in \mathcal{M}$ and

$$h = Af_m = Bf_m \in B(\mathcal{M})$$

Thus B is semi-Fredholm and $\text{ind } B = \dim \ker B$. Then there is a $\delta_1 > 0$ s.t. for $|\mu| < \delta_1$,

$$\dim \ker B - \mu \leq \dim \ker B, \quad \dim \text{ran } (B - \mu)^\perp = 0 \quad \text{and} \quad \text{ind } (B - \mu) = \text{ind } B$$

And thus $\dim \ker B - \mu = \dim \ker B$.

Also There is a δ_2 s.t. $\text{ind } (A - \mu) = \text{ind } A$ for all $|\mu| < \delta_2$. Then let $\delta = \min \{\delta_1, \delta_2\}$.

Claim: $\ker (A - \mu) \subset \mathcal{M}$ for $0 < |\mu| < \delta$.

If $h \in \ker (A - \mu)$, then $A^n h = \mu^n h$. So $h = A^n \mu^{-n} h \in \mathcal{M}$.

Finally, for $0 < |\mu| < \delta$,

$$\dim \ker (A - \mu) = \dim \ker (B - \mu) = \dim \ker B$$

And since $\text{ind } (A - \mu)$ is constant, $\dim \text{ran } (A - \mu)^\perp$ is also constant. \square

Theorem 10. *If $\lambda \in \partial\sigma(A)$, then either λ is an isolated point of $\sigma(A)$ or $\lambda \in \sigma_{le}(A) \cap \sigma_{re}(A)$.*

Proof. Suppose $\lambda \in \partial\sigma(A)$ and $\lambda \notin \sigma_{le}(A) \cap \sigma_{re}(A)$. Therefore, $A - \lambda$ is semi-Fredholm. By above theorem, there is a $\delta > 0$ s.t. $\dim \ker (A - \mu)$ and $\dim \text{ran } (A - \mu)^\perp$ are constant for $0 < |\lambda - \mu| < \delta$. Since $\lambda \in \partial\sigma(A)$, there is a $\nu \in \rho(A)$ with $0 < |\nu - \mu| < \delta$. Therefore,

$$\ker (A - \mu) = \{0\} = \text{ran } (A - \mu)^\perp, \quad \text{for } 0 < |\lambda - \mu| < \delta$$

And $A - \mu$ is semi-Fredholm, $\text{ran } (A - \mu)$ is closed, thus $A - \mu$ is invertible for $0 < |\lambda - \mu| < \delta$, which means that λ is isolated. \square

Proposition 11. *Let λ be an isolated point of $\sigma(A)$, then the following statements are equivalent.*

- 1) $\lambda \notin \sigma_{le}(A) \cap \sigma_{re}(A)$.

2) $A - \lambda$ is Fredholm and $\text{ind}(A - \lambda) = 0$.

3) The Riesz idempotent $E(\{\lambda\})$ is finite rank.

Proof. 1) \Rightarrow 2): Since $\lambda \notin \sigma_{le}(A) \cap \sigma_{re}(A)$, $\text{ran}(A - \lambda)$ is closed and either $\dim \ker(A - \lambda) < \infty$ or $\dim \ker(A - \lambda)^* < \infty$. Assume $\dim \ker(A - \lambda) < \infty$. Since λ is isolated, we can use similar argument in the **Theorem 8** in the subsection **5.1.2** to get

$$\dim \ker(A - \lambda) = \dim \ker(A - \lambda)^* < \infty$$

And if we assume $\dim \ker(A - \lambda)^* < \infty$, we can also get same result. Thus $A - \lambda$ is Fredholm and $\text{ind}(A - \lambda) = 0$.

2) \Rightarrow 3): By the definition,

$$(A - \lambda)E(\{\lambda\}) = \int_{\Gamma} \frac{(z - \lambda)\chi_{\{\lambda\}}(z)}{z - A} dz = 0 \Rightarrow E(\{\lambda\})\mathcal{H} \subset \ker(A - \lambda)$$

Thus $E(\{\lambda\})$ is finite rank.

3) \Rightarrow 1): As above mention, $E(\{\lambda\})\mathcal{H} \subset \ker(A - \lambda)$. In fact, the converse is also true.

Claim: $\ker(A - \lambda) \subset E(\{\lambda\})\mathcal{H}$.

Let $\Delta = \sigma(A) \setminus \{\lambda\}$, $\mathcal{H}_{\Delta} = E(\{\Delta\})\mathcal{H}$ and $A_{\Delta} = A|_{\mathcal{H}_{\Delta}}$, and similarly for $\{\lambda\}$. Then $A_{\Delta} - \lambda$ is invertible. If $h \in \ker(A - \lambda)$,

$$\begin{aligned} 0 = (A - \lambda)h &= (A - \lambda)E(\{\lambda\})h + (A - \lambda)E(\Delta)h \\ &= (A_{\lambda} - \lambda)E(\{\lambda\})h + (A_{\Delta} - \lambda)E(\Delta)h \end{aligned}$$

Since $E(\{\lambda\}) \perp E(\Delta)$,

$$(A_{\lambda} - \lambda)E(\{\lambda\})h = (A_{\Delta} - \lambda)E(\Delta)h = 0$$

And the fact that $A_{\Delta} - \lambda$ is invertible implies that $E(\Delta)h = 0$, thus $h \in \mathcal{H}_{\{\lambda\}}$.

Thus, $\ker(A - \lambda) = E(\{\lambda\})\mathcal{H}$ is finite. Moreover, by the similar argument in the **Theorem 8** in the subsection **5.1.2**, $\text{ran}(A - \lambda)$ is closed. Thus, $A - \lambda$ is semi-Fredholm. \square

5.5 Compact Perturbations

5.5.1 Weyl-von Neumann Theorem

In finite dimensional case, any normal operator can be diagonalized, but it is not true for the infinite dimensional case. In fact, it can hold by doing some extra efforts. Firstly, we can do that for a self-adjoint operator.

Like the definition of Hilbert-Schmidt operators, we define $B \in \mathcal{B}_p(\mathcal{H})$ if $|B|^{\frac{p}{2}} \in \mathcal{B}_2(\mathcal{H})$, and

$$\|B\|_p = \left\| |B|^{\frac{p}{2}} \right\|_2^{\frac{2}{p}} = \left(\sum_e \|B e\|^p \right)^{\frac{1}{p}}$$

Clearly, $\|B\| \leq \|B\|_p$. And moreover, if B has finite rank m , then

$$\|B\|_p \leq m^{\frac{1}{p}} \|T\|$$

In fact, $\mathcal{B}_p(\mathcal{H}) \subset \mathcal{B}_0(\mathcal{H})$.

Theorem 1 (Weyl-von Neumann Theorem). *If A is a self-adjoint operator on a separable Hilbert space, $\varepsilon > 0$ and $1 < p < \infty$, then there is a diagonalizable self-adjoint operator D s.t.*

$$A - D \in \mathcal{B}_p(\mathcal{H}), \text{ and } \|A - D\|_p < \varepsilon$$

To prove it, we need a lemma.

Lemma 2. *If A is a self-adjoint operator on a separable Hilbert space, $\varepsilon > 0$ and $1 < p < \infty$, then there is a finite rank projection P and a self-adjoint $K \in \mathcal{B}_p(\mathcal{H})$, s.t. $f \in P$, $\|K\|_p < \varepsilon$ and $A + K$ is reduced by P .*

Proof. Let E be the spectral measure for A and $\sigma(A) \subset [a, b]$ and intervals $\{\Delta_k\}_{k=1}^n$ be a partition of $[a, b]$ with length $(b - a)/n$. Put

$$f_k = E(\Delta_k)f, \text{ and } g_k = \frac{f_k}{\|f_k\|} \text{ for } f_k \neq 0$$

Let λ_k be the midpoint of Δ_k then

$$\|(A - \lambda_k)g_k\| \leq \frac{(b - a)}{n}$$

Define P as the projection onto $\text{span } g_1, \dots, g_n$. Since $P^\perp g_k = 0$,

$$\|P^\perp A g_k\| = \|P^\perp (A - \lambda_k)g_k\| \leq \frac{(b - a)}{n}$$

Since $g_k \in E(\Delta_k)$, $A g_k \in E(\Delta_k)$. Thus, $A g_k \perp g_j$ for $j \neq k$.

$$P^\perp A g_k = A g_k - \sum_{j=1}^n \langle A g_k, g_j \rangle g_j = A g_k - \langle A g_k, g_k \rangle g_k \in E(\Delta_k)$$

Thus $P^\perp Ag_k \perp P^\perp Ag_j$ for $j \neq k$. Then

$$\begin{aligned} \|P^\perp APh\| &= \left\| \sum_{k=1}^n \langle h, g_k \rangle P^\perp Ag_k \right\| \\ &= \sum_{k=1}^n |\langle h, g_k \rangle|^2 \|P^\perp Ag_k\|^2 \\ &\leq \left(\frac{(b-a)}{n} \right)^2 \|h\|^2 \end{aligned}$$

Therefore, $\|P^\perp AP\| \leq (b-a)/n$ and by above mention

$$\|P^\perp AP\|_p \leq n^{\frac{1}{p}} \left(\frac{(b-a)}{n} \right) = \frac{(b-a)}{n^q}, \quad q = 1 - \frac{1}{p}$$

Put $B = PAP + P^\perp AP^\perp$ and $K = -(P^\perp AP + PAP^\perp)$, then $B = A + K$. Clearly, K is self-adjoint and P reduces B and

$$\|K\|_p \leq \frac{2(b-a)}{n^q}$$

Thus for any $\varepsilon > 0$, $\|K\|_p < \varepsilon$ for some partition. \square

Then we can prove the Weyl-von Neumann Theorem.

Proof of Weyl-von Neumann Theorem. Let $\{h_j\}$ be a dense subset in \mathcal{H} and put $f = h_1$, then by above lemma, there is a finite rank projection P_1 and a self-adjoint operator K_1 s.t. $A + K_1$ is reduced by P_1 and $h_1 \in P_1$ and $\|K_1\|_p < \varepsilon/2$. Then applying above lemma to $(A + K_1)|_{P_1^\perp}$ and the vector $f = P_1^\perp h_2$, we get K_2 and P_2 and extend K_2 on \mathcal{H} by defining $K_2 h = 0$ for $h \in P_1$. By reduction, we get finite rank projections $\{P_n\}$ and $\{K_n\}$ s.t.

- 1) $P_j P_n = 0$ for $j \neq n$;
- 2) $\|K_n\|_p < \varepsilon/2^n$;
- 3) $h_n \in P_1 + \cdots + P_n$;
- 4) $A + K_1 + \cdots + K_n$ is reduced by $P_1 + \cdots + P_n$;
- 5) $K_n(P_1 + \cdots + P_{n-1}) = 0$.

Then let $K = \sum K_n$, then $\|K_n\|_p < \varepsilon$. And let $D = A + K$. Clearly, D is self-adjoint. Since $1 = \sum_n P_n$ and D is reduced by any P_n , if $D_n = D|_{P_n}$, $D = \oplus D_n$. And since P_n is finite, each D_n can be diagonalized, thus D is diagonalizable. \square

5.5.2 Weyl-von Neumann-Berg Theorem

Now, we can apply above theorem to the normal operators. But it can be more general. Since a class of commutative finite rank normal operators can be diagonalizable simultaneously, we want to see the similar result for the infinite dimensional case.

By the **Theorem 4** in the section 4.3, we can get a trivial corollary.

Corollary 1. *If N_1, \dots, N_n are commuting normal operators on a separable Hilbert space, then there is a self-adjoint operator A and continuous functions f_1, \dots, f_n s.t. $N_j = f_j(A)$ for $1 \leq j \leq n$.*

Theorem 2 (Weyl-von Neumann-Berg Theorem). *If N_1, \dots, N_n are commuting normal operators on a separable Hilbert space and $\varepsilon > 0$, then there are diagonalizable normal operators D_1, \dots, D_n and compact operators K_1, \dots, K_n s.t. for $1 \leq j \leq n$, $\|K_j\| < \varepsilon$ and $N_j = D_j + K_j$.*

By above corollary and the Weyl-von Neumann Theorem, $A = D + K$ with $\|K\| < \varepsilon$, but we need to prove that $f(A) - f(D) \in \mathcal{B}_0(\mathcal{H})$, and $\|f(A) - f(D)\| M < \varepsilon$.

Lemma 3. 1) *Let A be a self-adjoint operator with $A = D + K$ for some diagonalizable D and compact operator K . Then for any continuous function f on $\sigma(A)$, there is a compact operator K' s.t. $f(A) = f(D) + K'$.*

2) *If f is a continuous function on \mathbb{R} , then the map $A \rightarrow f(A)$ is uniformly continuous on a bounded sets of self-adjoint operators.*

Proof. For 1), if $A = D + K$, then for any $m \geq 1$, $A^m = D^m + K_m$ for some compact operator K_m . Therefore, for each polynomial p ,

$$p(A) = p(D) + K_p, \text{ for some operator } K_p \in \mathcal{B}_0(\mathcal{H})$$

Let $\{p_k\}$ be a sequence converging uniformly to f on a closed interval containing $\sigma(A) \cup \sigma(D)$, thus

$$\|p_k(A) - f(A)\| \rightarrow 0, \quad \|p_k(D) - f(D)\| \rightarrow 0$$

Therefore, we have

$$K_{pk} = p_k(A) - p_k(D) \rightarrow f(A) - f(D) = K' \text{ in norm}$$

and $K' \in \mathcal{B}_0(\mathcal{H})$.

For 2), we just need to prove this if f is a polynomial since on a compact set the limit of a sequence of uniformly continuous functions is also uniformly continuous. Note that

$$T^{n+1} - S^{n+1} = T^n(T - S) + (T^n - S^n)S$$

That means $T \rightarrow T^n$ is uniformly continuous on a bounded set of self-adjoint operators. Therefore, $A \rightarrow f(A)$ for any polynomial f . \square

Weyl-von Neumann-Berg Theorem. By above corollary, let A be the self-adjoint operator and f_1, \dots, f_n be continuous functions on $\sigma(A)$ s.t. $N_j = f_j(A)$ for any $1 \leq j \leq n$. By the Weyl-von Neumann Theorem, there are self-adjoint operators D and K with $K \in \mathcal{B}_0(\mathcal{H})$ and $\|K\| < \varepsilon$ s.t. $A = D + K$. Then by above lemma,

$$N_j - f_j(D) \in \mathcal{B}_0(\mathcal{H}), \|N_j - f_j(D)\| = \|f_j(A) - f_j(D)\| < \varepsilon \quad \square$$

Corollary 4. *If N is a normal operator on a separable Hilbert space and $\varepsilon > 0$, then there is a diagonalizable normal operator D s.t. $N = D \oplus K + K_1$ and $\sigma(D) = \sigma_e(D) = \sigma_e(N)$, where K and K_1 are compact.*

Proof. By above theorem, $N = D_1 + K_1$ for the normal diagonalizable operator D_1 and self-adjoint compact operator K_1 with $\|K_1\| < \varepsilon$. Then by the Fredholm Theory, $\sigma_e(D_1) = \sigma_e(N)$. By the **Theorem 5** in the subsection **5.4.3** and the separability of the Hilbert space,

$$\{\mu_j\}_{j=1}^{\infty} = \sigma(D_1) \setminus \sigma_e(D_1)$$

where each μ_j has finite multiplicity and is repeated as often as their multiplicity. Let $\{h_j\}$ be the corresponding eigenvectors then

$$D_1 = D \oplus \sum_{j=1}^{\infty} \mu_j h_j \otimes h_j$$

where D is a diagonalizable normal operator on a subspace of \mathcal{H} with $\sigma(D) = \sigma_e(D) = \sigma_e(D_1)$. Put

$$K = \sum_{j=1}^{\infty} \mu_j h_j \otimes h_j$$

then K is compact. Therefore,

$$N = D \oplus K + K_1 \quad \square$$

In above corollary, we cannot value the norm of K , but in some special case, we can do this. Firstly, we need some lemmas.

Lemma 5. *Let \mathfrak{A} be a Banach algebra and $a \in \mathfrak{A}$.*

- 1) *If $\alpha \in \rho(a)$, then $\text{dist}(\alpha, \sigma(a)) \geq \|(a - \alpha)^{-1}\|^{-1}$.*
- 2) *If $b, k \in \mathfrak{A}$ with $\|k\| < \varepsilon$ for some $\varepsilon > 0$ s.t. $a = b + k$, then*

$$\{ \alpha : \alpha \in \rho(a), \|(a - \alpha)^{-1}\|^{-1} > \varepsilon \} \subset \rho(b)$$

Proof. For 1), if $\alpha \in \rho(a)$ and by the perturbation of inverse,

$$|\beta| < \|(a - \alpha)^{-1}\|^{-1} \Rightarrow \beta + \alpha \in \rho(a)$$

Therefore, 1) holds.

For 2), if

$$\alpha \in \{ \alpha : \alpha \in \rho(a), \|(a - \alpha)^{-1}\|^{-1} > \varepsilon \}$$

then by the fact that $\|k\| < \varepsilon$,

$$\|(a - \alpha)^{-1}k\| \leq \|(a - \alpha)^{-1}\| \|k\| < 1$$

Therefore, $1 - (a - \alpha)^{-1}k$ is invertible. It shows

$$b - \alpha = a - \alpha - k = (a - \alpha)(1 - (a - \alpha)^{-1}k)$$

is invertible, i.e. $\alpha \in \rho(a)$. □

Lemma 6. *If $A \in \mathcal{B}(\mathcal{H})$ is a normal operator, then for any $\lambda \in \rho(A)$*

$$\|(A - \lambda)^{-1}\| = \text{dist}(\lambda, \sigma(A))^{-1}$$

Proof. Let E be the spectral measure for A , then

$$(A - \lambda)^{-1} = \int_{\sigma(A)} (z - \lambda)^{-1} dE(z)$$

Then, the norm

$$\|(A - \lambda)^{-1}\| = \sup \{ |(z - \lambda)^{-1}| : z \in \sigma(A) \} = \text{dist}(\lambda, \sigma(A))^{-1} \quad \square$$

Combining above lemmas, if $N = D + K$ with $\|K\| < \varepsilon$, then

$$\sigma(D) \subset \{ \lambda : \text{dist}(\lambda, \sigma(N)) \leq \varepsilon \}$$

Then we can refine above corollary.

Proposition 7. *If N is a normal operator on a separable Hilbert space with $\sigma(N) = \sigma_e(N)$ and $\varepsilon > 0$, then there is a diagonalizable normal operator D s.t. $N \cong D \oplus K + K_1$ and $\sigma(D) = \sigma_e(D) = \sigma_e(N)$, where K and K_1 are compact with $\|K\|, \|K_1\| < \varepsilon$*

Proof. By above corollary, $N = D_1 + K_1$ for the normal diagonalizable operator D_1 and self-adjoint compact operator K_1 with $\|K_1\| < \varepsilon$. And

$$D_1 = D_2 \oplus \sum_{j=1}^{\infty} \mu_j h_j \otimes h_j$$

where D_2 is a diagonalizable normal operator on a subspace of \mathcal{H} with $\sigma(D_2) = \sigma_e(D_2) = \sigma_e(D_1)$. Since

$$\sigma(D_1) = \sigma_e(D_1) \cup \{\mu_j\}, \quad \sigma_e(D_1) = \sigma_e(N) = \sigma(N)$$

then by above mention for any μ_j , there is a $\lambda_j \in \sigma_e(D_2)$, s.t.

$$\begin{aligned} |\mu_j - \lambda_j| &= \text{dist}(\mu_j, \sigma_e(D_2)) = \text{dist}(\mu_j, \sigma_e(D_1)) \\ &= \text{dist}(\mu_j, \sigma(N)) \\ &\leq \varepsilon \end{aligned}$$

Then since D_2 is diagonalizable, let g_j be the corresponding vector with respect to λ_j ,

$$D_1 = D \oplus \sum_{j=1}^{\infty} (-\lambda_j) g_j \otimes g_j \oplus \sum_{j=1}^{\infty} \mu_j h_j \otimes h_j$$

And in fact,

$$D_1 \cong D \oplus \sum_{j=1}^{\infty} (\mu_j - \lambda_j) h_j \otimes h_j$$

where D a diagonalizable normal operator on a subspace of \mathcal{H} with $\sigma(D) = \sigma_e(D) = \sigma_e(N)$. Put

$$K = \sum_{j=1}^{\infty} (\mu_j - \lambda_j) h_j \otimes h_j$$

then K is compact with $\|K\| < \varepsilon$. Therefore,

$$N \cong D \oplus K + K_1$$

□

5.5.3 Approximately Unitary Equivalence

For two operators N and M , if there is a unitary U s.t. $N = U^*MU$, N and M are same in some sense and we say that these two are unitary equivalent. But this relation is too strong, so there is another weaker relation, approximately unitary equivalence.

Lemma 1. *If N and M are diagonalizable normal operators on a separable Hilbert space with*

$$\sigma(N) = \sigma_e(N) = \sigma_e(M) = \sigma(M)$$

*and $\varepsilon > 0$, then there is a unitary operator U s.t. $N - U^*MU \in \mathcal{B}_1(\mathcal{H})$ with $\|N - U^*MU\|_1 < \varepsilon$.*

Proof. By the condition, if ν is an isolated point of $\sigma(N)$, it must be an eigenvalue with infinite multiplicity, and thus it is an eigenvalue of M with infinite multiplicity. If ν is a limit point of $\sigma(N)$, then there is a sequence in $\sigma(M)$ converging to ν .

Let $\{\nu_n\}$ and $\{\mu_m\}$ be all eigenvalues of N and M respectively and repeated as often as their multiplicity. Since $\sigma(N) = \sigma_e(N)$, all of them have infinite multiplicity. Moreover, for any ν_j , there is a sequence $\{\mu_{n_j}\}$ converging to it.

Claim: There is a renumbering of the two sets $\{\nu_{n_k}\}$ and $\{\mu_{m_k}\}$ s.t.

$$|\nu_{n_k} - \mu_{m_k}| < \frac{\varepsilon}{2^k}, \quad \forall k \in \{n_1, \dots, n_{2k}\} \cap \{m_1, \dots, m_{2k}\}$$

Step 1: Let $n_1 = 1$ and m_1 be the smallest integer with

$$|\nu_1 - \mu_{m_1}| < \frac{\varepsilon}{2}$$

Step 2: If $m_1 > 1$, put $m_2 = 1$, otherwise, put $m_2 = 2$. Then let ν_2 be the smallest integer larger than 1 s.t.

$$|\nu_2 - \mu_{m_2}| < \frac{\varepsilon}{2^2}$$

Step 3: Continuing Step 1 and Step 2 to get $\{n_k\}$ and $\{m_k\}$ for $1 \leq k \leq 2j$. Then let n_{2j+1} be the smallest integer and $n \neq n_k$ for $1 \leq k \leq 2j$ and m_{2j+1} be the smallest integer greater than m_k for $1 \leq k \leq 2j$ with

$$|\nu_{n_{2j+1}} - \mu_{m_{2j+1}}| < \frac{\varepsilon}{2^{2j+1}}$$

Similarly, we get n_{2j+2} and m_{2j+1} with same properties. Therefore, by induction, this claim holds.

Since N and M are diagonalizable on a separable Hilbert space, there are two orthonormal basis $\{e_n\}$ and $\{f_m\}$ s.t.

$$Ne_n = \nu_n e_n, \text{ and } Mf_m = \mu_m f_m$$

Define U as $Ue_{n_k} = f_{m_k}$ for all k , then clearly U is a unitary. And,

$$(N - U^*MU)e_{n_k} = (\nu_{n_k} - \mu_{m_k})e_{n_k}$$

thus $\|N - U^*MU\|_1 < \varepsilon$. □

Theorem 2. *If N and M are normal operators on a separable Hilbert space with same essential spectrum, then there is a unitary U s.t. $N - U^*MU$ is compact.*

Proof. According the last corollary in above subsection, there are diagonalizable normal operators N_1 and M_1 s.t.

$$N = N_1 \oplus K_1 + K, \quad M = M_1 \oplus K'_1 + K'$$

where K_1, K and K'_1, K' are compact operators. And since $\sigma_e(N) = \sigma_e(M)$,

$$\sigma(N_1) = \sigma_e(N_1) = \sigma_e(M_1) = \sigma(M_1)$$

By above lemma, there is a unitary V s.t. $N_1 - V^*M_1V \in \mathcal{B}_0(\mathcal{H})$. Therefore, put $U = V \oplus id$,

$$N - U^*MU = (N_1 - V^*M_1V) \oplus (K_1 - K'_1) + U^*(K - K')U$$

is compact. □

By above theorem, any two normal operators with same essential spectrum actually differ by some compact operator. By we cannot estimate the norm of this compact operator. However, the last proposition in above subsection can provide us a clue. If the two normal operators N and M satisfy

$$\sigma(N) = \sigma_e(N) = \sigma_e(M) = \sigma(M)$$

then there are similar N_1, K_1, K and M_1, K'_1, K' s.t.

$$N = N_1 \oplus K_1 + K, \quad M = M_1 \oplus K'_1 + K'$$

But the difference is that the norm of K_1 and K'_1 can be arbitrary small. Therefore, by above theorem, the unitary U can be choosing s.t. $N - U^*MU$ has arbitrary small norm.

In general case, we need a definition to describe this definition.

Definition 3. Two normal operators A and B are approximately unitarily equivalent, denoted by $A \cong_a B$, if there is a sequence of unitaries $\{U_n\}$ s.t.

$$\|A - U_n^*BU_n\| \rightarrow 0$$

Moreover, if $A - U_n^*BU_n \in \mathcal{B}_0(\mathcal{H})$ and $A \cong_a B$, then A and B are called strongly approximately unitarily equivalent, denoted by $A \cong_{sa} B$.

Lemma 4. If \mathfrak{A} is a C^* -algebra and $\{a_n\}$ is a sequence of normal elements of \mathfrak{A} s.t. $a_n \rightarrow a$ in norm, then $\sigma(a_n) \rightarrow \sigma(a)$.

Proof. Since \mathfrak{A} is a C^* -algebra, by the continuity of the involution and the multiplication, a is normal. By the Continuous Functional Calculus, we know

$$C^*(a_n) \cong C(\sigma(a_n)), \quad C^*(a) \cong C(\sigma(a))$$

For any $f_k \in C(\sigma(a_k))$, extending f_k on $\cup\sigma(a_n) \cup \sigma(a)$ by setting $f_k = 0$ on other areal. Then $f_k \in B(\cup\sigma(a_n) \cup \sigma(a))$, and $\|f_k\|$ does not change. And this mention can also apply to any f in $C(\sigma(a))$.

Let $g_k(z) = z \in C(\sigma(a_k))$ and $g(z) = z \in C(\sigma(a))$, then extending g_k and g on $\cup\sigma(a_n) \cup \sigma(a)$. Therefore, $\|g_k\| = \|a_k\|$ and $\|g\| = \|a\|$. Take $\alpha \in \sigma(a)$, then since $a_n \rightarrow a$ in norm, $g_n \rightarrow g$ in norm (?), therefore for any $0 < \varepsilon < |g(\alpha)|$, there is a n_0 s.t. for all $n > n_0$

$$|g_n(\alpha) - g(\alpha)| \leq \|g_n - g\| < \varepsilon$$

Then we can choose some $\alpha_n \in g_n^{-1}(g_n(\alpha)) \in \sigma(\alpha_n)$ since $g_n(\alpha) \neq 0$. Therefore, $\alpha_n \rightarrow \alpha$. And the converse is similar. \square

Theorem 5. *If N and M are normal operators on a separable Hilbert spaces, then the following statements are equivalent.*

- 1) $N \cong_{sa} M$.
- 2) $N \cong_a M$.
- 3) $\sigma_e(N) = \sigma_e(M)$ and for any $\lambda \in \sigma_e(N) = \sigma_e(M)$,

$$\dim \ker (N - \lambda) = \dim \ker (M - \lambda)$$

Proof. Let E_N and E_M be the spectral measures for N and M . 1) implies 2) clearly.

2) \Rightarrow 3): Let $\{U_n\}$ be the unitaries s.t. $\|N - U_n^* M U_n\| \rightarrow 0$. Then by above lemma,

$$X = \sigma(M) = \sigma(N)$$

If $\lambda \in X$ is isolated, then λ is eigenvalue for M and N . And the characteristic function $\chi_{\{\lambda\}}(z)$ satisfies

$$\chi_{\{\lambda\}}(U_n^* M U_n) \longrightarrow \chi_{\{\lambda\}}(N)$$

in norm. That means $E_N(\{\lambda\})$ and $E_M(\{\lambda\})$ has same rank, thus

$$\dim \ker (N - \lambda) = \dim \ker (M - \lambda)$$

for any isolated point λ in X . Since the essential spectrum consists of all points in the spectrum that are not isolated eigenvalues of finite multiplicity, thus $\sigma_e(M) = \sigma_e(N)$.

3) \Rightarrow 1): By the Spectral Theorem,

$$N = N_1 \oplus N_2, \text{ and } M = M_1 \oplus M_2,$$

where $\sigma(M_2) = \sigma_e(M_2) = \sigma_e(M) = \sigma_e(N) = \sigma_e(N_2) = \sigma(N_2)$ and the spectrum of N_1 and M_1 consists of all eigenvalues of N and M with finite

multiplicity, and thus $N_1 \cong M_1$. Let V be the unitary s.t. $V^*M_1V = N_1$. Moreover, since

$$\sigma(M_2) = \sigma_e(M_2) = \sigma_e(N_2) = \sigma(N_2)$$

by above mention there is a sequence of unitaries $\{W_n\}$ s.t.

$$W_n^*M_2W_n \longrightarrow N_2$$

in norm and $N_2 - W_n^*M_2W_n \in \mathcal{B}_0(\mathcal{H})$. Then putting $U_n = W_n \oplus V_n$, thus $N \cong_{sa} M$. \square

Chapter 6

Von Neumann Algebras

Abelian von Neuman algebras can be explicitly described as $L^\infty(\mu)$. For general von Neumann algebras, they have more complex structure but more interesting properties. Firstly, we have seen the relation between the *WOT* and the ultraweak topology, therefore, the ultraweak topology can provide some different things of the von Neumann algebra. Then, since the weak closedness of von Neumann algebras, all projections in a von Neumann algebra can completely generate the whole algebra. Thus projections are fundamental elements in the von Neumann algebra. Researching the properties of them and classifying them can help us to classify general von Neumann algebras. By the properties of the projections contained in a von Neumann algebras, general von Neumann algebras can be classified as 5 different categories and we can see that the abelian one is belong to what category.

6.1 General Von Neumann Algebras

6.1.1 Elementary Properties

Proposition 1. *If $\{A_i\}$ is an increasing net of self-adjoint operators on \mathcal{H} s.t. $\sup_i \|A_i\| < \infty$, then there is an operator A s.t. the following hold.*

- 1) $A_i \leq A$ for all i and if B is any self-adjoint operator with $A_i \leq B$ for all i , then $A \leq B$. Thus $A = \sup_i A_i$.
- 2) $A_i \rightarrow A$ in *WOT*.
- 3) $A_i \rightarrow A$ in *SOT*.
- 4) $A_i \rightarrow A$ in *weak**.

Proof. Without loss of generality, assume all A_i are positive.

For any $h \in \mathcal{H}$, $\{\langle A_i h, h \rangle\}$ is an increasing positive sequence and bounded by $\alpha \|h\|^2$ for $\alpha = \sup_i \|A_i\|$. Then

$$F(h, h) = \lim_i \langle A_i h, h \rangle, \quad |F(h, h)| \leq \alpha \|h\|^2$$

Clearly, F is a sesquilinear and by the Polar Identity,

$$|F(g, h)| \leq \alpha \|g\| \|h\|$$

Then there is a $A \in \mathcal{B}(\mathcal{H})$ s.t.

$$\langle Ag, h \rangle = F(g, h) = \lim_i \langle A_i g, h \rangle$$

Therefore, A is positive, $A \rightarrow A_i$ in WOT , $A_i \leq A$ and $\|A\| \leq \alpha$. Then

$$0 \leq A - A_i \leq A \Rightarrow (A - A_i)^{\frac{1}{2}} \leq A^{\frac{1}{2}}, \forall i$$

If $h \in \mathcal{H}$, then

$$\begin{aligned} \|(A - A_i)h\|^2 &\leq \left\| (A - A_i)^{\frac{1}{2}} \right\|^2 \left\| (A - A_i)^{\frac{1}{2}} h \right\|^2 \\ &\leq \|A\| \left\| (A - A_i)^{\frac{1}{2}} h \right\|^2 \\ &= \|A\| \langle (A - A_i)h, h \rangle \rightarrow 0 \end{aligned}$$

Thus $A_i \rightarrow A$ in SOT .

Finally, to see that $A_i \rightarrow A$ in $weak^*$, consider $\{A_i^{(\infty)}\}$, and by 2), there is a $A^{(\infty)}$ s.t. $A_i^{(\infty)} \rightarrow A^{(\infty)}$ in WOT , therefore, $A_i \rightarrow A$ in $weak^*$. \square

Corollary 2. *If \mathfrak{A} is a C^* -algebra contained in $\mathcal{B}(\mathcal{H})$ but without the identity and \mathfrak{A} is $weak^*$ -closed, then there is a projection P in \mathfrak{A} s.t.*

$$A = PA = AP, \forall A \in \mathfrak{A}$$

Therefore, P is the identity for \mathfrak{A} and $\mathfrak{A}|_P$ is a von Neumann algebra.

Proof. There is an approximate identity $\{A_i\}$. By above proposition,

$$P = \sup_i A_i \in \mathfrak{A}, \text{ and } A_i \rightarrow P \text{ in } weak^*$$

Then for any $A \in \mathfrak{A}$, $\|AA_i - A\| \rightarrow 0$ and $AA_i \rightarrow AP$ in $weak^*$. Thus $A = AP = PA$. \square

Proposition 3. 1) *If \mathfrak{A} is a von Neumann algebra and \mathcal{E} is a family of projections in \mathfrak{A} , then $\bigvee \mathcal{E}$ and $\bigwedge \mathcal{E}$ belong to \mathfrak{A} .*

2) *If \mathfrak{A} is a von Neumann algebra and $1 \leq d \leq \infty$, then $\mathfrak{A}^{(d)}$ is a von Neumann algebra and*

$$(\mathfrak{A}^{(d)})' = \{ T = (T_{ij}) \in \mathcal{B}(\mathcal{H}^{(d)}): T_{ij} \in \mathfrak{A}' \}$$

3) If \mathfrak{A}_i is a von Neumann algebra in $\mathcal{B}(\mathcal{H}_i)$ for all i and $\mathcal{H} = \bigoplus_i \mathcal{H}_i$, then

$$\mathfrak{A} = \bigoplus_i \mathfrak{A}_i = \{ \bigoplus_i A_i : A_i \in \mathfrak{A}_i, \sup_i \|A_i\| < \infty \}$$

is a von Neumann algebra.

Proposition 4. Let $\{T_i\}$ and $\{S_i\}$ be nets in $\mathcal{B}(\mathcal{H})$.

- 1) If $T_i \rightarrow 0$ in WOT and $S_i \rightarrow$ in SOT , then $S_i T_i, T_i S_i \rightarrow 0$ in SOT .
- 2) If $T_i \rightarrow 0$ in SOT and $K \in \mathcal{B}_0(\mathcal{H})$, then $T_i K$ and $K T_i$ convergence to 0 in norm. If $T_i \rightarrow 0$ in WOT then $T_i K \rightarrow 0$ in norm and $K T_i \rightarrow 0$ in SOT .

Proof. For 1), since $T_i \rightarrow 0$ in WOT , by the Principle of Uniform Boundedness, $\{T_i\}$ is norm boundedness. Since for any $h \in \mathcal{H}$, $\|S_i h\| \rightarrow 0$,

$$\|T_i S_i h\| \leq \|T_i\| \|S_i h\| \rightarrow 0$$

And $\|S_i T_i h\| \rightarrow 0$ clearly. To prove 2), firstly for finite rank K_n and let h_1, \dots, h_n and g_1, \dots, g_n be in \mathcal{H} ,

$$K_n = \sum_{i=1}^n h_i \otimes g_i$$

That means $\text{ran } K_n = \text{span}\{h_1, \dots, h_n\}$. Let P_n be the projection onto $\text{ran } K_n$, then $T_i K_n = T_i P_n$,

$$T_i K_n = T_i P_n : \text{ran } K_n \longrightarrow \text{span}\{T_i h_1, \dots, T_i h_n\}$$

And since the WOT and the SOT and the norm topology are same on a finite dimensional space, $T_i K_n \rightarrow 0$ in norm for $T_i \rightarrow 0$ in WOT or in SOT . For any $K \in \mathcal{B}_0(\mathcal{H})$, there is a sequence of finite rank operators $\{K_n\}$ s.t. $K_n \rightarrow K$ in norm. Since

$$\|T_i K\| = \|T_i K - T_i K_n\| + \|T_i K_n\|$$

$T_i K$ convergences to 0 in norm for $T_i \rightarrow 0$ in WOT or in SOT .

$$K_n T_i = \sum_{l=1}^n h_n \otimes (T_i^* g_n)$$

Thus,

$$\begin{aligned} \|K_n T_i\| &\leq \sum_{l=1}^n \|h_n \otimes (T_i^* g_n)\| \\ &= \sum_{l=1}^n \|h_n\| \|T_i^* g_n\| \end{aligned}$$

If $T_i \rightarrow 0$ in SOT , then $\|T_i^* g_n\| \rightarrow 0$ for all n , i.e. $\|K_n T_i\| \rightarrow 0$. Thus for $K \in \mathcal{B}_0(\mathcal{H})$, $KT_i \rightarrow 0$ in norm. If $T_i \rightarrow 0$ in WOT , for any $x \in \mathcal{H}$,

$$\begin{aligned} \langle K_n T_i x, K_n T_i x \rangle &= \left\langle \sum_{l=1}^n \langle T_i x, h_l \rangle g_l, \sum_{m=1}^n \langle T_i x, h_m \rangle g_m \right\rangle \\ &= \sum_{l=1}^n \sum_{m=1}^n \langle T_i x, h_l \rangle \langle T_i x, h_m \rangle \langle g_l, g_m \rangle \end{aligned}$$

Therefore, if $T_i \rightarrow 0$ in WOT , $\|KT_i x\| \rightarrow 0$ for $x \in \mathcal{H}$, i.e. $KT_i \rightarrow 0$ in SOT . \square

6.1.2 Central Cover

Definition 1. For a von Neumann algebra \mathfrak{A} , if $\mathcal{Z} = \mathfrak{A} \cap \mathfrak{A}'$ is the center of \mathfrak{A} , the central cover or central cover of $A \in \mathfrak{A}$ is the projection

$$C_A = \inf \{ C : C \in \mathcal{Z} \text{ is a projection, } AC = A \}$$

Remark. By above proposition, C_A is definitely a projection satisfying

$$AC_A = C_A = A$$

Proposition 2. Let \mathfrak{A} be a von Neumann algebra.

1) If $A \in \mathfrak{A}$, $E = (\ker A)^\perp$ and $F = \overline{\text{ran } A}$, then

$$C_A = C_{A^*} = C_{A^*A} = C_{AA^*} = C_E = C_F$$

2) If E is a projection in \mathfrak{A} , then

$$C_E = \overline{\text{span} \{ Ah : h \in E, A \in \mathfrak{A} \}}$$

Proof. For 1), $C_A = C_{A^*}$ is trivial by the definition. If C is a central projection s.t. $AC = A$, then $A^*AC = A^*A$. Conversely, if $A^*AC = A^*A$, then $|A|C = |A|$ by taking the root. By the Polar Decomposition, $AC = A$. Thus $C_A = C_{A^*A}$.

Firstly, $E \in \mathfrak{A}$. If C is a central projection s.t. $AC = A$, then

$$C(\ker A) \subset \ker A \Rightarrow EC = CE = E \Rightarrow E \leq C$$

Therefore, $C_E \leq C_A$. Conversely, if C is a central projection s.t. $E \leq C$, then

$$AC = AEC = AE = A$$

Thus $C_A \leq C_E$. Moreover, $C_F = C_E$.

For 2), let

$$P = \overline{\text{span}\{Ah: h \in E, A \in \mathfrak{A}\}}$$

and clearly $P \in \mathfrak{A}'$. Moreover, for any $B \in \mathfrak{A}'$ and $h \in E$,

$$BAh = ABh \in P \Rightarrow BP \subset P \Rightarrow P \in \mathfrak{A}'' = \mathfrak{A}$$

Therefore, P is a central projection with $E \leq P$, thus $C_E \leq P$. Conversely, if C is a central projection s.t. $EC = E$, then for $A \in \mathfrak{A}$ and $h \in E$,

$$C(Ah) = AC(Eh) = AEh = Ah \Rightarrow CP = P = PC$$

Thus $P \leq C$, then $P \leq C_E$. □

Definition 3. If \mathfrak{A} is a von Neumann algebra and E is a projection, defining

$$\mathfrak{A}_E = \{ A_E: A \in \mathfrak{A} \}$$

where $A_E = EA|_E$.

Remark. In general, \mathfrak{A}_E is not an algebra. However, if $E \in \mathfrak{A}$, there is a $*$ -isomorphism between \mathfrak{A}_E and $E\mathfrak{A}E$. If $E \in \mathfrak{A}'$, there is a $*$ -homomorphism between \mathfrak{A}_E and $E\mathfrak{A}E$. $E\mathfrak{A}E$ is clear a C^* -subalgebra of \mathfrak{A} .

Proposition 4. Let \mathfrak{A} be a von Neumann algebra.

- 1) If $E \in \mathfrak{A}'$ is a projection, then \mathfrak{A}_E is a von Neumann algebra in $\mathcal{B}(E)$ and $(\mathfrak{A}_E)' = (\mathfrak{A})'_E$.
- 2) If $E \in \mathfrak{A}$ is a projection, then \mathfrak{A}_E is a von Neumann algebra in $\mathcal{B}(E)$ and $(\mathfrak{A}_E)' = (\mathfrak{A})'_E$.
- 3) If \mathcal{Z} is the center of \mathfrak{A} and E is a projection in \mathfrak{A} or \mathfrak{A}' , then the center of \mathfrak{A}_E is \mathcal{Z}_E .

Proof. For 1), it is clear since E reduces \mathfrak{A} .

Assume that $E \in \mathfrak{A}$. \mathfrak{A}_E and $(\mathfrak{A}')_E$ are $*$ -algebra clearly. If $B \in \mathfrak{A}'$ and $A \in \mathfrak{A}$, then for $h \in E$,

$$B_E A_E h = EBEAh = EAE Bh = A_E B_E h$$

Therefore, $(\mathfrak{A})'_E \subset (\mathfrak{A}_E)'$.

Check: $(\mathfrak{A}_E)' \subset (\mathfrak{A})'_E$

Let $U \in (\mathfrak{A}_E)' \subset \mathcal{B}(E)$ be a unitary, $A_1, \dots, A_n \in \mathfrak{A}$ and $h_1, \dots, h_n \in E$,

then

$$\begin{aligned}
\left\| \sum_{k=1}^n A_k U h_k \right\|^2 &= \sum_{i,j}^n \langle A_i U h_i, A_j U h_j \rangle = \sum_{i,j}^n \langle (E A_j^* A_i E) U h_i, U h_j \rangle \\
&= \sum_{i,j}^n \langle U (E A_j^* A_i E) h_i, U h_j \rangle = \sum_{i,j}^n \langle E A_j^* A_i E h_i, h_j \rangle \\
&= \left\| \sum_{k=1}^n A_k h_k \right\|^2
\end{aligned}$$

Therefore, for $C = \overline{\text{span}\{Ah : h \in E, A \in \mathfrak{A}\}}$, there is a partial isometry with the initial space C ,

$$W\left(\sum_{k=1}^n A_k h_k\right) = \sum_{k=1}^n A_k U h_k$$

and $WC^\perp = 0$. By above proposition $C = C_E$. Therefore, for any $A \in \mathfrak{A}$,

$$WAC^\perp = WC^\perp A = 0 = AW C^\perp$$

and on C ,

$$WA \sum_{k=1}^n A_k h_k = W \sum_{k=1}^n A A_k h_k = A \sum_{k=1}^n A_k U h_k = AW \sum_{k=1}^n A_k h_k$$

Therefore, $W \in \mathfrak{A}'$ with $Wh = Uh$ for $h \in E$. So $W_E = U$. Since any element in $(\mathfrak{A}_E)'$ can be the linear combination of four unitaries, $(\mathfrak{A}_E)' \subset (\mathfrak{A})'_E$.

Check: $\mathfrak{A}_E \subset \mathcal{B}(E)$ is von Neumann.

Let $T \in (\mathfrak{A}_E)''$ and extending T on \mathcal{H} by setting $T = 0$ on E^\perp . If $B \in \mathfrak{A}'$, then $B_E \in (\mathfrak{A}_E)'$. For $h \in E$,

$$TBh = (TE)Bh = TB_E h = B_E Th = B_E Th = BTh$$

And on the other hand,

$$TBE^\perp = TE^\perp B = 0 = BTE^\perp$$

Thus $T \in \mathfrak{A}'' = \mathfrak{A}$, and $(\mathfrak{A}_E)'' = \mathfrak{A}_E$.

For 3), by above mention,

$$\mathfrak{A}_E \cap (\mathfrak{A}_E)' = \mathfrak{A}_E \cap (\mathfrak{A})'_E = \mathcal{Z}_E$$

□

6.1.3 Kaplansky Density Theorem

For C^* -algebra \mathfrak{A} , let $\mathcal{B}(\mathfrak{A})$ be the norm-closed unit ball in \mathfrak{A} .

Theorem 1 (Kaplansky Density Theorem). *If \mathfrak{B} is a C^* -subalgebra of $\mathcal{B}(\mathcal{H})$ and $\mathfrak{A} = \overline{\mathfrak{B}}^{SOT}$, then*

- 1) $\mathcal{B}(\mathfrak{A})$ is the SOT closure of $\mathcal{B}(\mathfrak{B})$.
- 2) $\mathcal{B}(\text{Re } \mathfrak{A})$ is the SOT closure of $\mathcal{B}(\text{Re } \mathfrak{B})$.
- 3) $\mathcal{B}(\mathfrak{A})_+$ is the SOT closure of $\mathcal{B}(\mathfrak{B})_+$.
- 4) the SOT closure of the unitaries in \mathfrak{B} contains the unitaries of \mathfrak{A} .

Lemma 2. *If $f \in C(\mathbb{R})$ with $f(0) = 0$ and there are positive constants a, b s.t.*

$$|f(t)| \leq a|t| + b, \quad \forall t \in \mathbb{R}$$

then for any net $\{T_i\} \subset \text{Re } \mathfrak{A}$ with $T_i \rightarrow T$ in SOT, $f(T_i) \rightarrow f(T)$ in SOT.

Proof. For any $g \in C(\mathbb{R})$, then the map $g: \text{Re } \mathcal{B}(\mathcal{H}) \rightarrow \text{Re } \mathcal{B}(\mathcal{H})$ makes sense. Denote the set of all such g s.t. g is continuous w.r.t. the SOT by \mathcal{S} and let \mathcal{S}_b be set of all norm bounded elements in \mathcal{S} . Clearly, $\mathcal{S}_b \mathcal{S} \subset \mathcal{S}$.

Let $e(t) = t/(1+t^2)$ and $A, B \in \text{Re } \mathcal{B}(\mathcal{H})$,

$$\begin{aligned} e(A) - e(B) &= (1+A^2)^{-1}(A(1+B^2) - (1+A^2)B)(1+B^2)^{-1} \\ &= (1+A^2)^{-1}(A-B)(1+B^2)^{-1} \\ &\quad + (1+A^2)^{-1}A(B-A)B(1+B^2)^{-1} \end{aligned}$$

Then for any $h \in \mathcal{H}$,

$$\|e(A) - e(B)\| \leq \|(A-B)(1+B^2)^{-1}h\| + \|(A-B)(1+B^2)^{-1}h\|$$

Therefore, $e \in \mathcal{S}_b$. Put $e_\alpha(t) = e(\alpha t)$, then $e_\alpha(t) \in \mathcal{S}_b$ for each $\alpha > 0$. Since $\{e_\alpha: \alpha > 0\}$ can separate the points in $\mathbb{R} \setminus \{0\}$ and \mathcal{S}_b is closed,

$$C_0(\mathbb{R} \setminus \{0\}) \subset \mathcal{S}_b$$

If f satisfy above assumption, $t \rightarrow f(t)/(1+t^2)$ belongs to $C_0(\mathbb{R} \setminus \{0\})$, and hence to \mathcal{S}_b . Since the identity function also belongs to \mathcal{S} and the condition f satisfies,

$$t \rightarrow \frac{tf(t)}{1+t^2} \in \mathcal{S}_b$$

And thus

$$t \rightarrow \frac{t^2 f(t)}{1+t^2} \in \mathcal{S}$$

Therefore,

$$f(t) = \frac{t^2 f(t)}{1+t^2} + \frac{f(t)}{1+t^2} \in \mathcal{S}$$

□

Then by using above lemma and constructing some special functions, we can easily prove the Kaplansky Theorem.

proof of Kaplansky Theorem. Fix a $A \in \mathcal{B}(\text{Re } \mathfrak{A})$ and a net $\{B_i\} \subset \text{Re } \mathfrak{B}$ s.t. $B_i \rightarrow A$ in SOT .

For 2), if $f = (t \wedge 1) \vee (-1)$, then f satisfies above lemma. Therefore,

$$f(B_i) \longrightarrow f(A) = A \text{ in } SOT$$

but $f(B_i) \in \mathcal{B}(\text{Re } \mathfrak{B})$.

Similarly, 3) holds by constructing the function $f(t) = (t \wedge 1) \vee 0$.

Let $U \in \mathfrak{A}$ be a unitary. Then by the functional calculus, there is a $A \in \text{Re } \mathfrak{A}$ s.t.

$$U = \exp iA = \cos A + i \sin A$$

Then let $\{B_i\} \subset \text{Re } \mathfrak{B}$ be a net s.t. $B_i \rightarrow A$ in SOT . Since both $\sin t$ and $\cos t - 1$ satisfy above lemma,

$$U_i = \cos B_i + i \sin B_i \longrightarrow U \text{ in } SOT$$

and $U_i \in \mathfrak{B}$ is a unitary.

For 1), fix a $A \in \mathcal{B}(\mathfrak{A})$ and consider the C^* -algebras $M_2(\mathfrak{B})$ and $M_2(\mathfrak{A})$, then $M_2(\mathfrak{A})$ is the SOT -closure of $M_2(\mathfrak{B})$.

$$\begin{pmatrix} 0 & A \\ A^* & 0 \end{pmatrix} \in \mathcal{B}(\text{Re } M_2(\mathfrak{A}))$$

therefore, by 1) there is a net $\{T_i\} \subset \mathcal{B}(\text{Re } M_2(\mathfrak{B}))$ converges to above element in SOT . Put

$$T_i = \begin{pmatrix} X_i & B_i \\ B_i^* & Y_i \end{pmatrix}$$

Since $\|T_i\| \leq 1$, $\|B_i\| \leq 1$. Clearly, $B_i \rightarrow A$ in SOT . □

6.2 Topologies on Von Neumann Algebras

6.2.1 Pedersen Up-Down Theorem

The order of all self-adjoint operators in von Neumann algebras can provide another important property of von Neumann algebras.

Definition 1. For any subset $\mathcal{S} \subset \text{Re } \mathcal{B}(\mathcal{H})$, define

$$\begin{aligned} \mathcal{S}^\sigma &= \{ T : \exists \{S_n\} \subset \mathcal{S} \text{ increasing sequence, } S_n \rightarrow T \text{ in } SOT \} \\ \mathcal{S}_\sigma &= \{ T : \exists \{S_n\} \subset \mathcal{S} \text{ decreasing sequence, } S_n \rightarrow T \text{ in } SOT \} \end{aligned}$$

Remark. Note that $\mathcal{S} \subset \mathcal{S}^\sigma$ and $\mathcal{S}_\sigma = -(-\mathcal{S})^\sigma$.

Lemma 2. *If A and B are two positive operators with $A \leq B$, then*

$$A(1 + A)^{-1} \leq B(1 + B)^{-1}$$

Proof. Because of $A(1 + A)^{-1} = 1 - (1 + A)^{-1}$ and the fact if $0 \leq A \leq B$, $B^{-1} \leq A^{-1}$, the lemma holds. \square

Lemma 3. *If \mathfrak{B} is a C^* -algebra and $\mathfrak{A} = \overline{\mathfrak{B}}^{SOT}$ and P is a projection in \mathfrak{A} , then for each sequence of unit $\{h_n\}$ in \mathcal{H} there is a $A \in ((\mathcal{B}(\mathfrak{B}_+))^\sigma)_\sigma$ s.t.*

$$A(1 - P)h_n = 0 = (1 - A)Ph_n, \quad \forall n \geq 1$$

Proof. By the Kaplansky Density Theorem and the sparability of \mathcal{H} , there is a sequence $\{C_n\}$ in $\mathcal{B}(\mathfrak{B}_+)$ s.t. $C_n \rightarrow P$ in SOT . Then it can assume that for $1 \leq j \leq n$,

$$\|C_n(1 - P)h_j\| < \frac{1}{n2^n}, \quad \|(1 - C_n)Ph_j\| < \frac{1}{n}$$

For $n < m$ define

$$A_{nm} = \left(1 + \sum_{k=n}^m kC_k\right)^{-1} \sum_{k=n}^m kC_k$$

Then each $A_{nm} \in \mathcal{B}(\mathfrak{B}_+)$ and $A_{nm} \leq \sum_{k=n}^m kC_k$. Thus for $j \leq n$,

$$\langle A_{nm}(1 - P)h_j, (1 - P)h_j \rangle \leq \sum_{k=n}^m \frac{1}{2^k} < \frac{1}{2^{n-1}}$$

Since $\sum_{k=n}^m kC_k \geq mC_m$, by above lemma

$$A_{nm} \leq (1 + mC_m)^{-1}(mC_m) \Rightarrow 1 - A_{nm} \leq (1 + mC_m)^{-1}$$

By the fact for $0 \leq t \leq 1$, $(1 + mt)^{-1} \leq (1 + m)^{-1}(1 + m(1 - t))$ and $C_m \in \mathcal{B}(\mathfrak{B}_+)$,

$$1 - A_{nm} \leq (1 + m)^{-1}(1 + m(1 - C_m))$$

Then it implies that

$$\begin{aligned} \langle (1 - A_{nm})Ph_j, Ph_j \rangle &\leq \frac{1}{m} \langle (1 + m(1 - C_m))Ph_j, Ph_j \rangle \\ &\leq \frac{1}{m} (\|Ph_j\| + m \|(1 - C_m)Ph_j\|) \\ &\leq \frac{2}{m} \end{aligned}$$

Fix a n , $\{A_{nm} : m > n\}$ is increasing by above lemma. Since it is bounded, there is a A_n s.t. $A_{nm} \rightarrow A_n$ fin SOT and $A_n \in (\mathcal{B}(\mathfrak{B}_+))^\sigma$. On the other

hand, when $n+1 < m$, also by above lemma, it implies that $A_{n+1,m} \leq A_{nm}$, thus $A_{n+1} \leq A_n$. Thus there is a operator A s.t. $A_n \rightarrow A$ in SOT , and thus $A \in ((\mathcal{B}(\mathfrak{B}_+))^\sigma)_\sigma$. Moreover, by above

$$\begin{aligned} \langle A_n(1-P)h_j, (1-P)h_j \rangle &\leq \frac{1}{2^{n-1}} \\ \langle (1-A_n)Ph_j, Ph_j \rangle &\leq 0 \end{aligned}$$

As $n \rightarrow \infty$, $A(1-P)h_n = 0 = (1-A)Ph_n$ for all n . \square

Theorem 4 (Pedersen Up-Down Theorem). *If \mathcal{H} is a separable, \mathfrak{B} is a C^* -algebra contained in $\mathcal{B}(\mathcal{H})$, and $\mathfrak{A} = \overline{\mathfrak{B}}^{SOT}$, then*

$$\mathcal{B}(\mathfrak{A}_+) = ((\mathcal{B}(\mathfrak{B}_+))^\sigma)_\sigma \text{ and } \text{Re } \mathfrak{A} = (\text{Re } \mathfrak{B})^\sigma_\sigma$$

Proof. Firstly, by above lemma, for any projection $P \in \mathfrak{A}$ and a dense subset in $\{h \in \mathcal{H} : \|h\| = 1\}$, there is a $A \in ((\mathcal{B}(\mathfrak{B}_+))^\sigma)_\sigma$ s.t. $A(1-P) = 0 = (1-A)P$.

For $A \in (\mathcal{B}(\mathfrak{A}_+))$ with the spectral measure E and $k \leq 1$,

$$P_k = E \left(\bigcup_{j=1}^{2^k-1} \left(\frac{j}{2^k}, \frac{j+1}{2^k} \right] \right)$$

then by the integral, $A = \sum_{k=1}^{\infty} 2^{-k} P_k$ converging by the norm. Then for each k , there is a decreasing sequence $\{Z_{kn}\}$ in $(\mathcal{B}(\mathfrak{B}_+))^\sigma$ s.t. $Z_{kn} \rightarrow P_k$ in SOT , define

$$T_n = \sum_{k=1}^n \frac{1}{2^k} Z_{kn} + \frac{1}{2^n} \geq A$$

Since $(\mathcal{B}(\mathfrak{B}_+))^\sigma$ is convex, $T_n \in (\mathcal{B}(\mathfrak{B}_+))^\sigma$. Moreover, because of $Z_{+1,n+1} < 1$,

$$T_n - T_{n+1} = \sum_{k=1}^n \frac{1}{2^k} (Z_{kn} - Z_{k,n+1}) + \frac{1}{2^n} \left(\frac{1}{2^{n+1}} Z_{n+1,n+1} + \frac{1}{2^{n+1}} \right) \leq 0$$

Thus $\{T_n\}$ is decreasing and there is some $T \in ((\mathcal{B}(\mathfrak{B}_+))^\sigma)_\sigma$ with $T \geq A$ s.t. $T_n \rightarrow T$ in SOT .

Check: $A = T \in ((\mathcal{B}(\mathfrak{B}_+))^\sigma)_\sigma$

Since $0 \leq Z_{kn} - P_k \leq 1$, the series $\sum_k 2^{-k} (Z_{kn} - P_k)$ is convergent. Then for any unit $h \in \mathcal{H}$ and $\varepsilon > 0$, chose N s.t. $2^{-N} < \varepsilon$, then there is an n_0 s.t.

$$0 \leq \langle (Z_{kn} - P_k)h, h \rangle < \frac{\varepsilon}{2}, \text{ for } 1 \leq k \leq N \text{ and } n \geq n_0$$

Therefore,

$$\lim_{n \rightarrow \infty} \sum_{k=1}^{\infty} \frac{1}{2^k} \langle (Z_{kn} - P_k)h, h \rangle = 0, \forall h \in \mathcal{H}$$

However,

$$0 \leq T_n - A = T_n - \sum_{k=1}^{\infty} \frac{1}{2^k} P_k \leq \frac{1}{2^n} + \sum_{k=1}^{\infty} \frac{1}{2^k} (Z_{kn} - P_k)$$

Therefore, $T_n \rightarrow A$ in WOT , but $T_n \rightarrow T$ in SOT , thus $T = A$.

Check: $\text{Re } \mathfrak{A} = (\text{Re } \mathfrak{A})^\sigma_\sigma$

Fix $A \in \text{Re } \mathfrak{A}$ and let $\alpha = -\|A\|$. Therefore, $A + \alpha \leq 0$. Let $\beta = \|A + \alpha\|$, so $T = \beta^{-1}(A + \alpha) \in \mathcal{B}(\mathfrak{A}_+)$ and $A = \beta T - \alpha$. But $\beta T \in ((\mathcal{B}(\mathfrak{B}_+))^\sigma)_\sigma$ and $-\alpha \in \mathcal{B}(\mathfrak{B})$. Thus $A \in ((\mathcal{B}(\text{Re } \mathfrak{B}))^\sigma)_\sigma$. \square

The Up-Down Theorem can reveal the basic property of general von Neumann algebras.

Theorem 5. *If \mathfrak{A} is a C^* -algebra contained in $\mathcal{B}(\mathcal{H})$, then \mathfrak{A} is weak*-closed if and only if it contains the supremum of every bounded increasing net of self-adjoint operators in the algebra.*

Proof. \mathfrak{A} is weak*-closed, thus \mathfrak{A} is a von Neumann algebra. Therefore, by above proposition \mathfrak{A} has this property. Assume that \mathfrak{A} is a C^* -algebra and if $\{A_i\}$ is an increasing net in $\text{Re } \mathfrak{A}$ with $A_i \rightarrow A$ in SOT , then $A \in \mathfrak{A}$. Then by taking the negatives of the increasing parts, it is also true for that if $\{A_i\}$ is an decreasing net in $\text{Re } \mathfrak{A}$ with $A_i \rightarrow A$ in SOT , then $A \in \mathfrak{A}$. Therefore,

$$\text{Re } \mathfrak{A} = ((\mathcal{B}(\text{Re } \mathfrak{A}))^\sigma)_\sigma$$

Then if \mathcal{H} is separable, by the Up-Down Theorem \mathfrak{A} is weak*-closed.

For the general case, let $\mathfrak{B} = \overline{\mathfrak{A}}^{SOT}$. It is sufficient to show that any projection contained in \mathfrak{B} is contained in \mathfrak{A} . Let P be the projection in \mathfrak{B} and $g \in P$ and $h \in P^\perp$. Then by above lemma, there is an $T \in \mathfrak{A}_+$ s.t. $Tg = g$ and $Th = 0$. Let R be projection onto $\overline{\text{ran } T}$ and $Rg = g$ and $Rh = 0$, moreover $R \in \mathfrak{A}$ because of the hypothesis and the Spectral Theorem. For any $g \in P$ and $h \in P^\perp$, denote the R by R_{gh} . Fix $g \in P$ and

$$\mathcal{F} = \{h_1, h_2, \dots, h_m\} \subset P^\perp$$

then $R_{\mathcal{F}} = R_{gh_1} \wedge R_{gh_2} \wedge \dots \wedge R_{gh_m} \in \mathfrak{A}$. And $\{R_{\mathcal{F}}\}$ is a decreasing net and thus $R_g = \lim -SOT R_{\mathcal{F}} \in \mathfrak{A}$. And since $R_g P^\perp = 0$, $R_g \leq P$ for all $g \in P$. Then let

$$\mathcal{G} = \{g_1, g_2, \dots, g_n\} \subset P$$

$R_{\mathcal{G}} = R_{g_1} \vee R_{g_2} \vee \dots \vee R_{g_n} \in \mathfrak{A}$. And $\{R_{\mathcal{G}}\}$ is an increasing net and thus $R = \lim -SOT R_{\mathcal{G}} \in \mathfrak{A}$. Therefore, $R = P \in \mathfrak{A}$. \square

6.2.2 Normal Homomorphisms

Definition 1. If \mathfrak{A} and \mathfrak{B} are von Neumann algebras and $\rho: \mathfrak{A} \rightarrow \mathfrak{B}$ is a linear positive map, then ρ is normal if for any increasing net $\{A_i\}$ in \mathfrak{A} with $A_i \rightarrow A$ in SOT , $\rho(A_i) \rightarrow \rho(A)$ in SOT .

Remark. In particular, if $E \in \mathfrak{A}$, $A \mapsto EAE$ is a normal homomorphism. The normal homomorphisms should be positive linear maps that are continuous with respect to the SOT , but in above definition, there is an extra condition. In fact, this extra condition has an equivalent expression.

If \mathfrak{A} is a von Neumann algebra, the \mathfrak{A} is $weak^*$ -closed, therefore \mathfrak{A} is the dual space of some Banach space. Let \mathfrak{A}_* be the space of all $weak^*$ -continuous functionals on \mathfrak{A} , then $\mathfrak{A} = (\mathfrak{A}_*)^*$. And also

$$\mathfrak{A}_* \cong \mathcal{B}_1(\mathcal{H})/\mathfrak{A}_\perp,$$

where \mathfrak{A}_\perp is the space of all trace class annihilating the \mathfrak{A} . Therefore, for any $L \in \mathfrak{A}_*$, there is $T \in \mathcal{B}_1(\mathcal{H})$ s.t. for all $A \in \mathfrak{A}$,

$$L(A) = \text{tr}(AT)$$

In fact, if L is positive, then T can be chosen as a positive element.

Lemma 2. If ϕ and ψ are positive functionals on a von Neumann algebra \mathfrak{A} s.t. there is an operator $A \in \mathfrak{A}_+$ with $\phi(A) < \psi(A)$, then there is an operator $B \in \mathfrak{A}_+$ s.t. $B \leq A$ and $\phi(T) < \psi(T)$ for all $T \in \mathfrak{A}$ with $0 < T \leq B$. If A is projection, then B can be chosen as a projection.

Proof. Let \mathcal{C} be the set of all $C \in \mathfrak{A}_+$ s.t. $C \leq A$ and $\phi(C) \leq \psi(C)$. And give \mathcal{C} with the usual order. Then if $\{C_i\}$ is the chain in \mathcal{C} and let $C = \sup_i C_i \in \mathfrak{A}_+$ and $C \leq A$. Since ϕ and ψ is positive, $C \in \mathcal{C}$. Then by the Zorn's Lemma, there is a maximal element C in \mathcal{C} . Put $B = A - C$. If $T \in \mathfrak{A}_+$ with $0 < T \leq B$, by the maximality of C , $T \notin \mathcal{C}$, thus $\phi(T) < \psi(T)$. Moreover, if A is a projection, $B \leq A$ implies that $\text{ran } B \leq A$. If $B = \int_0^1 t dP(t)$ be the spectral decomposition, then for any $\varepsilon > 0$, $P = P([\varepsilon, 1])$. Then P satisfies above condition. \square

Theorem 3. If ψ is a positive linear functional on the von Neumann algebra \mathfrak{A} , then the following statements are equivalent.

- 1) ψ is normal.
- 2) If $\{E_i\}$ is a pairwise orthogonal family of projections in \mathfrak{A} , then

$$\psi\left(\sum E_i\right) = \sum \psi(E_i)$$

- 3) ψ is $weak^*$ -continuous.

4) *There is a positive trace class operator C s.t. $\psi(A) = \text{tr}(AC)$ for all $A \in \mathfrak{A}$.*

Proof. Without loss the generality, assuming that $\psi(1) = 1$.

1) \Rightarrow 2): It is trivial.

2) \Rightarrow 3): Since $0 \leq \psi(E) \leq 1$ for any non-zero projection E , for any $h \in \mathcal{H}$,

$$\psi(E) \leq \langle Eh, h \rangle$$

Therefore, by preceding lemma, there is a projection F with $0 \leq F \leq E$ s.t.

$$\psi(T) \leq \langle hT, h \rangle, \quad \forall 0 \leq T \leq F$$

Since $FAF \leq \|A\| F$,

$$\psi(FAF) \leq \|A\| \langle FAFh, h \rangle$$

Then by the CBS inequality, for any $A \in \mathcal{B}(\mathfrak{A})$

$$|\psi(AF)|^2 \leq \psi(FA^*AF) \leq \|A\|^2 \langle FA^*AFh, h \rangle = \|A\|^2 \|AFh\|^2$$

Therefore, $\psi(\cdot F) : \mathfrak{A} \rightarrow \mathbb{C}$ is *SOT*-continuous.

Let $\{E_i\}$ be the maximal family of pairwise orthogonal projections in \mathfrak{A} s.t. $\psi(\cdot E_i)$ is *SOT*-continuous on \mathfrak{A} . Let $E = \sum_i E_i$. Clearly, $E = 1$. If $E \neq 1$, then there is a projection $F \leq E^\perp$ s.t. $\psi(\cdot F)$ is *SOT*-continuous on \mathfrak{A} , which contradicts the maximality of $\{E_i\}$. Then by 2),

$$\psi\left(\sum_i E_i\right) = \psi(1) = 1$$

Therefore, for any $\varepsilon > 0$, there is a finite set I_0 of index s.t. for any J with $I_0 \subset J$, and let $P_J = \sum_J E_j$,

$$\psi P_J^\perp = 1 - \psi(P_J) < \varepsilon$$

Then for any $A \in \mathcal{B}(\mathfrak{A})$,

$$\left| \psi(AP_J^\perp) \right|^2 \leq \psi(AA^*)\psi(P_J) < \varepsilon$$

And thus for any J with $I_0 \subset J$,

$$\|\psi - \psi(\cdot P_J)\| = \sup \{ |\psi(AP_J)| : A \in \mathcal{B}(\mathfrak{A}) \} < \sqrt{\varepsilon}$$

Since $\psi(\cdot P_J) \in \mathfrak{A}_*$, $\psi \in \mathfrak{A}_*$.

3) \Rightarrow 4): Firstly, assume that $\psi(A) = \langle Ag, h \rangle$ for some $g, h \in \mathcal{H}$.

Check: $\psi(A) = \langle Af, f \rangle$ for some $f \in \mathcal{H}$.

By the fact that $A \in \mathfrak{A}_+$, $\langle Ag, h \rangle \leq 0$,

$$\begin{aligned} 4\langle Ag, h \rangle &= \langle A(g+h), g+h \rangle - \langle A(g-h), g-h \rangle \\ &\leq \langle A(g+h), g+h \rangle \end{aligned}$$

Let $\phi(A) = \frac{1}{4}\langle A(g+h), g+h \rangle$, then $\psi(A) \leq \phi(A)$. By the **Proposition 7** in the subsection **3.2.7**, there is a $T \in \mathfrak{A}'$ with $0 \leq T \leq 1$ s.t.

$$\psi(A) = \langle AT(g+h), g+h \rangle$$

Let $f = T^{\frac{1}{2}}(g+h)$, then $\psi(A) = \langle Af, f \rangle = \text{tr}(Af \otimes f)$, where $f \otimes f$ is a positive trace class.

For general case, since $\psi \in \mathfrak{A}_*$, by above mention, there is a $D \in \mathcal{B}_0(\mathcal{H})$ s.t. $\psi(A) = \text{tr}(AD)$. By the **Theorem 3** in the subsection **5.3.3**, there are $g, h \in \mathcal{H}^{(\infty)}$ s.t.

$$\psi(A) = \text{tr}(AD) = \langle A^{(\infty)}g, h \rangle = \langle A^{(\infty)}f, f \rangle$$

Then there is a $C \in \mathcal{B}_0(\mathcal{H})$, s.t.

$$\langle A^{(\infty)}f, f \rangle = \text{tr}(AC)$$

And in fact, if $f = (f_i) \in \mathcal{H}^{(\infty)}$, $C = \sum_i f_i \otimes f_i \geq 0$.

4) \Rightarrow 1): It is trivial by definition. \square

Then by using the above theorem, the equivalent statement as the following corollary shows can make the definition of normal homomorphism more nature.

Corollary 4. *If \mathfrak{A} and \mathfrak{B} are von Neumann algebras and $\rho: \mathfrak{A} \rightarrow \mathfrak{B}$ is a positive linear map, then ρ is normal if and only if it is weak*-continuous.*

Proof. If ρ is weak*-continuous, then it is clearly normal. Conversely, if ρ is normal and $\phi \in \mathfrak{B}_*$ such that ϕ is positive, then $\phi \circ \rho$ is a normal functional on \mathfrak{A} . By preceding theorem and the fact that all element in \mathfrak{B}_* can be linear combination of four positive elements, $\phi \circ \rho$ is weak*-continuous for all $\phi \in \mathfrak{B}_*$. Therefore, by the Hahn-Banach Theorem, ρ is weak*-continuous. \square

Proposition 5. *Every *-isomorphism between von Neumann algebras is normal.*

Proof. Let $\rho: \mathfrak{A} \rightarrow \mathfrak{B}$ be *-isomorphism between von Neumann algebras \mathfrak{A} and \mathfrak{B} . If $\{\mathfrak{A}_i\}$ is an increasing net of self-adjoint operators in \mathfrak{A} with

$A = \sup_i A_i$, then $\{\rho(\mathfrak{A}_i)\}$ is an increasing net in \mathfrak{B} with $\rho(\mathfrak{A}_i) \leq \rho(A)$ for all A_i . Thus

$$B = \sup_i \rho(A_i) \leq \rho(A)$$

Conversely, $A_i = \rho^{-1}\rho(A_i) \leq \rho^{-1}(B)$ for all A_i . That means $A \leq \rho^{-1}(B)$. Therefore, $A = \rho(B)$. \square

6.2.3 Ideals

Rather than considering norm closed ideals in C^* -algebras, in von Neumann algebras, WOT -closed ideal should be paid more attention on, or equivalently, $weak^*$ -closed ideal.

Theorem 1. *Let \mathfrak{A} be a von Neumann algebra.*

- 1) \mathfrak{B} is a $weak^*$ -closed hereditary subalgebra of \mathfrak{A} if and only if there is a unique projection P in \mathfrak{A} s.t. $\mathfrak{B} = P\mathfrak{A}P$.
- 2) \mathfrak{J} is a $weak^*$ -closed left ideal of \mathfrak{A} if and only if there is a unique projection P in \mathfrak{A} s.t. $\mathfrak{J} = \mathfrak{A}P$.
- 3) \mathfrak{J} is a $weak^*$ -closed ideal of \mathfrak{A} if and only if there is a unique central projection P in \mathfrak{A} s.t. $\mathfrak{J} = \mathfrak{A}P = P\mathfrak{A}$.

Proof. For 1), if $\mathfrak{B} = P\mathfrak{A}P = \mathfrak{A}P \cap P\mathfrak{A} = \mathfrak{A}P \cap (\mathfrak{A}P)^*$, \mathfrak{B} is hereditary by the fact that $\mathfrak{A}P$ is a norm closed left ideal and **Theorem 3** in the subsection 3.2.5. And clearly, \mathfrak{B} is $weak^*$ -closed subalgebra. Conversely, if \mathfrak{B} is a $weak^*$ -closed hereditary subalgebra, then by the **Corollary 2** in the subsection 6.1.1, there is a central projection $P \in \mathfrak{B}$ s.t. $\mathfrak{B} = P\mathfrak{B} = \mathfrak{B}P$, thus $\mathfrak{B} \subset P\mathfrak{A}P$. If $A \in \mathfrak{A}_+$, then $PAP \leq \|A\|P \in \mathfrak{B}$. Since \mathfrak{B} is hereditary, $PAP \in \mathfrak{B}$. Therefore, $\mathfrak{B} = P\mathfrak{A}P$. The uniqueness of P is because any such P is the identity for \mathfrak{B} .

For 2), if \mathfrak{J} is a $weak^*$ -closed left ideal of \mathfrak{A} , then $\mathfrak{B} = \mathfrak{J} \cap \mathfrak{J}^*$ is a $weak^*$ -closed hereditary subalgebra. Then there is a projection $P \in \mathfrak{B} \subset \mathfrak{J}$ s.t. $\mathfrak{B} = P\mathfrak{A}P$. Thus $\mathfrak{A}P \subset \mathfrak{J}$. Conversely, if $T \in \mathfrak{J}$, then $T^*T \in \mathfrak{B}_+ = \mathfrak{J}_+$. Therefore, $T^*TP = T^*T$ and so $|T|P = |T|$. By the Polar Decomposition, $T = TP \in \mathfrak{A}P$. Hence, $\mathfrak{J} = \mathfrak{A}P$. And the uniqueness of P is because the uniqueness of P in \mathfrak{B} .

For 3), That \mathfrak{J} is a $weak^*$ -closed ideal of \mathfrak{A} , also norm closed, implies that $\mathfrak{J} = I^*$. Then by 1) and 2), there is a projection $P \in \mathfrak{J}$ s.t. $\mathfrak{J} = \mathfrak{A}P = P\mathfrak{A}$. Moreover, for any $A \in \mathfrak{A}$,

$$PA = P(PA) = P(PAP) = (PA)P = AP$$

Therefore, P is in the center of \mathfrak{A} . \square

For the ideal in a von Neumann algebra, it is always assumed as the *weak**-closed ideal. Thus, there is a direct corollary of above theorem about the simple von Neumann algebras.

Corollary 2. *A von Neumann algebra \mathfrak{A} is simple if and only if*

$$\mathcal{Z} = \mathfrak{A} \cap \mathfrak{A}' = \mathbb{C}I$$

Definition 3. A von Neumann algebra \mathfrak{A} is called factor, if

$$\mathcal{Z} = \mathfrak{A} \cap \mathfrak{A}' = \mathbb{C}I$$

Remark. By above corollary, factors are simple von Neumann algebras. Moreover, for a von Neumann algebra \mathfrak{A} , if $P \in \mathcal{Z}$ is a projection, then $\mathfrak{A} = P\mathfrak{A} + P^\perp\mathfrak{A}$, where $P\mathfrak{A}$ and $P^\perp\mathfrak{A}$ are von Neumann algebras. And by above mention, $P\mathfrak{A} = P\mathfrak{A}P \cong \mathfrak{A}_P$. Therefore, any von Neumann algebra is semi-simple and since $\mathfrak{I} = P\mathfrak{A}$ is an ideal, $\mathfrak{A}/\mathfrak{I} \cong P^\perp\mathfrak{A}$

Theorem 4. *If \mathfrak{A} and \mathfrak{B} are von Neumann algebras and $\rho: \mathfrak{A} \rightarrow \mathfrak{B}$ is a normal $*$ -homomorphism, then $\ker \rho$ is *weak**-closed and $\rho(\mathfrak{A})$ is *weak**-closed in \mathfrak{B} . Moreover, there is a central projection P s.t. $\ker \rho = P^\perp\mathfrak{A}$ and ρ is a $*$ -isomorphism of $P\mathfrak{A}$ onto $\text{ran } \rho$.*

Proof. Check: $\ker \rho$ is *weak**-closed.

If $\{A_i\}$ is a bounded increasing net in $\text{Re } \ker \rho$, then put $A = \sup_i A_i$. Since ρ is normal, $\rho(A) = \sup_i \rho(A_i) = 0$. Therefore, $A \in \ker \rho$. Then by **Theorem 5** in the subsection **6.2.1**, $\ker \rho$ is *weak**-closed.

Check: $\text{ran } \rho$ is *weak**-closed.

Let $\{B_i\}$ is a bounded increasing net in $\text{Re } \text{ran } \rho$ and $B = \sup_i B_i$. Then put

$$\Gamma = \{ A \in \mathfrak{A}_+ : \rho(A) \in \{B_i\} \}$$

Claim: Γ is directed with respect to the usual order.

Suppose $A_i, A_j \in \Gamma$ with $\rho(A_i) = B_i$ and $\rho(A_j) = B_j$. Pick B_k s.t. $B_i, B_j \leq B_k$. Therefore, there is a $T \in \mathfrak{A}_+$ s.t. $\rho(T) = B_k - B_i$. If $C_i = A_i + T$, then $C_i \geq A_i$ and $\rho(C_i) = B_k$. Similarly, there is a $C_j \geq A_j$ s.t. $\rho(C_j) = B_k$. Therefore, $C_i - C_j \in \ker \rho$, and if $A = C_i + |C_i - C_j|$, then $\rho(A) = B_k$ and $C_i, C_j \leq A$.

Fix $\varepsilon > 0$, by proof of the Pedersen Up-Down Theorem, the map $A \rightarrow (1 + \varepsilon A)^{-1}A$ is an increasing function on $\mathcal{B}(\mathcal{H})_+$. Thus $\{(1 + \varepsilon A)^{-1}A : A \in \Gamma\}$ is an increasing function in \mathfrak{A} with the bound ε^{-1} . Put

$$A_\varepsilon = \sup \{(1 + \varepsilon A)^{-1}A : A \in \Gamma\}$$

and $\|A_\varepsilon\| \leq \varepsilon^{-1}$. Since ρ is normal,

$$\rho(A_\varepsilon) = (1 + \varepsilon B)^{-1}B \in \text{ran } \rho$$

But $(1 + \varepsilon B)^{-1}B \rightarrow B$ in norm and $\text{ran } \rho$ is norm closed, thus $B \in \text{ran } \rho$.

The rest of the theorem can be obtained by using above theorem. \square

6.3 Projections

6.3.1 Equivalence of Projections

Projections in a von Neumann algebra play the role as generators, which lead to researching global properties of the von Neumann algebra is based on the properties of its projections, to some extent. Firstly, there is some equivalent relation between projections.

Definition 1. Let E and F be two projections in a von Neumann algebra \mathfrak{A} . E and F are equivalent, denoted by $E \sim F$, if there is a partial isometry $W \in \mathfrak{A}$ s.t. $W^*W = E$ and $WW^* = F$. If there is a $G \leq F$ s.t. $E \sim F$, then denoting $E \lesssim F$.

There are some properties of the equivalent projections.

Proposition 2. 1) In an abelian algebra, two projections are equivalent if and only if they are equal.

2) Two projections in $\mathcal{B}(\mathcal{H})$ are equivalent if and only if they have same dimensional.

3) If \mathfrak{A} is a von Neumann algebra and $A \in \mathfrak{A}$, then $\overline{\text{ran}(A)} \sim \overline{\text{ran}(A^*)}$.

4) Let E be a projection in a von Neumann algebra \mathfrak{A} . $E \sim 0$ if and only if $E = 0$.

5) If E and F are two projections in \mathfrak{A} with $E \sim F$, then for any central projection Z , $ZE \sim ZF$.

Proof. 1) holds trivially. For 2), by the proof of the Polar Decomposition, for a partial isometry W , $(\ker W)^\perp$ and $\text{ran } W$ are isometrically isomorphism. For 3), let $A = W|A|$ the Polar Decomposition. Then

$$W^*W = \overline{\text{ran}(|A|)} = \overline{\text{ran}(A^*)}, \quad WW^* = \overline{\text{ran}(A)}$$

For 4), $E = W^*W$ and $WW^* = 0$. That means $W = 0$, thus $E = 0$.

5) holds since $Z = Z^*Z$ and Z is in center. \square

Lemma 3. For projections E and F :

$$1) \overline{\text{ran}(EF)} = E \wedge (E \wedge F^\perp)^\perp = E - E \wedge F^\perp;$$

$$2) \overline{\text{ran}(EF^\perp)} = E - E \wedge F;$$

$$3) \overline{\text{ran}(F^\perp E)} = E \vee F - F.$$

Proof. It is sufficient to prove 1).

Check: $\ker FE = E^\perp + (E \vee F^\perp)$.

If $h \in \ker FE$ and expressing $h = g + f$ where $g \in E^\perp$ and $f \in E$, then $0 = FEh = Ff$. Thus $f \in F^\perp$, i.e. $f \in E \vee F^\perp$. Therefore,

$$\ker FE \leq E^\perp + (E \vee F^\perp)$$

Conversely, if $h = g + f$ with $g \in E^\perp$ and $f \in E \vee F^\perp$, then $FEh = Ff = 0$. Thus $\ker FE = E^\perp + (E \vee F^\perp)$.

Thus,

$$\overline{\text{ran}(EF)} = (\ker FE)^\perp = E \wedge (E \vee F^\perp)^\perp = E - E \vee F^\perp \quad \square$$

Proposition 4. 1) If E and F are projections, then

$$\begin{aligned} E \vee F - F &\sim E - E \wedge F \\ E - E \wedge F^\perp &\sim F - E^\perp \wedge F \end{aligned}$$

2) If $\{E_i\}$ and $\{F_i\}$ are two sets of pairwise orthogonal projections s.t. $E_i \sim_{W_i} F_i$ for all i , then

$$E = \sum_i E_i \sim_W F = \sum_i F_i$$

where $W = \sum_i W_i$.

Proof. For 1), by above proposition,

$$\overline{\text{ran}(F^\perp E)} \sim \overline{\text{ran}(F^\perp E^*)} = \overline{\text{ran}(EF^\perp)}$$

And by above lemma,

$$\overline{\text{ran}(F^\perp E)} = E \vee F - F, \text{ and } \overline{\text{ran}(EF^\perp)} = E - E \wedge F$$

Therefore, the first equivalence is obtained. By $\overline{\text{ran}(FE)} \sim \overline{\text{ran}(EF)}$, the second is also true.

For 2), if $h \in \mathcal{H}$,

$$\begin{aligned} \|Eh\|^2 &= \langle Eh, h \rangle = \sum_i \langle E_i h, h \rangle \\ &= \sum_i \langle W_i^* W_i h, h \rangle = \sum_i \|W_i\|^2 \end{aligned}$$

Since $\{W_i h\}$ is an orthogonal set, $\sum_i W_i$ converges strongly. Therefore put $W = \sum_i W_i$. And by the fact that $\{E_i\}$ are pairwise orthogonal, for any finite subset I of index,

$$\left(\sum_{i \in I} W_i\right)^* \left(\sum_{i \in I} W_i\right) = \sum_{i \in I} W_i^* W_i = \sum_{i \in I} E_i$$

Then by SOT convergence, $\sum_i W_i^* W_i = \sum_i E_i$, and it is also valid for F . \square

Proposition 5. *If E and F are two projections in \mathfrak{A} , then the following statements are equivalent.*

- 1) *There are nonzero projections E_1 and F_1 in \mathfrak{A} s.t. $E_1 \leq E$ and $F_1 \leq F$, and $E_1 \sim F_1$*
- 2) $C_E C_F \neq 0$.
- 3) $E\mathfrak{A}F \neq \{0\}$.

Proof. 1) \Rightarrow 2): By above lemma, $C_{E_1} = C_{F_1}$ and $C_{E_1} \leq C_E$ and $C_{F_1} \leq C_F$, thus $C_E C_F \neq 0$.

2) \Rightarrow 3): Since $C_E = \overline{\text{span } \mathfrak{A}E}$ and $C_F = \overline{\text{span } \mathfrak{A}F}$, if $E\mathfrak{A}F = \{0\}$, $EC_F = 0$. Thus $C_E \leq C_F^{bot}$, which is a contradiction.

3) \Rightarrow 2): If $A \in \mathfrak{A}$ s.t. $EAF \neq 0$, let $E_1 = \overline{\text{ran}(EAF)}$ and $F_1 = \overline{\text{ran}(EAF^*)} = \overline{\text{ran}(FA^*E)}$. Therefore, $E_1 \leq E$ and $F_1 \leq F$, and $E_1 \sim F_1$. \square

Theorem 6 (Comparison Theorem). *If E and F are two projections in \mathfrak{A} , then there is a central projection Z s.t.*

$$ZE \lesssim ZF, \text{ and } Z^\perp F \lesssim Z^\perp E$$

Proof. Finding Z .

By Zorn's Lemma, there is a maximal family $\{(E_i, F_i)\}$ of pairs of projections in \mathfrak{A} s.t. for all i ,

$$E_i \sim F_i, E_i \leq E, F_i \leq F$$

and $\{E_i\}$ and $\{F_i\}$ are pairwise orthogonal. Then by above proposition,

$$E' = \sum_i E_i \sim \sum_i F_i = F'$$

Then put $E'' = E - E'$ and $F'' = F - F'$. Therefore by the maximality and above proposition, $C_{E''} C_{F''} = 0$, then put $Z = C_{F''}$.

Check: $ZE \lesssim ZF$ and $Z^\perp F \lesssim Z^\perp E$.

Let W be the partial isometry s.t. $E' = W^*W$ and $F' = WW^*$. Then $ZE' \sim_{ZW} ZF'$. Therefore,

$$EZ = E'Z + E''Z = E'Z + E''C_{E''}C_{F''} = E'Z \sim ZF' \leq ZF$$

Hence, $ZE \lesssim ZF$. Similarly, $Z^\perp F \lesssim Z^\perp E$. \square

Corollary 7. *If \mathfrak{A} is a factor and E and F are two projections in \mathfrak{A} , then either $E \lesssim F$ or $F \lesssim E$.*

Lemma 8. *If \mathcal{L} is a complete lattice and $\tau: \mathcal{L} \rightarrow \mathcal{L}$ is an order preserving map, then there is an $x_0 \in \mathcal{L}$ s.t. $x_0 = \tau(x_0)$.*

Proof. Put $T = \{x \in \mathcal{L} : x \leq \tau(x)\}$ and $x_0 = \sup T \in \mathcal{L}$. If $x \in T$, then

$$x \leq x_0, \text{ and } x \leq \tau(x) \leq \tau(x_0) \Rightarrow x_0 \leq \tau(x_0)$$

Conversely, above inequality implies that $\tau(x_0) \leq \tau(\tau(x_0))$, thus $\tau(x_0) \in T$. Therefore, $x_0 \geq \tau(x_0)$. \square

Theorem 9. *If E and F are two projections in \mathfrak{A} s.t. $E \lesssim F$ and $F \lesssim E$, then $E \sim F$.*

Proof. Let W and V and two partial isometry s.t. $E = W^*W$ and $WW^* \leq F$ and $F = V^*V$ and $VV^* \leq E$. Let \mathcal{L} be the set of all projection G in \mathfrak{A} s.t. $G \leq F$. Then \mathcal{L} is complete lattice. Define the map τ on \mathcal{L} as

$$\begin{aligned} \tau: \quad \mathcal{L} &\longrightarrow \mathcal{L} \\ G &\longmapsto F - W^*(E - V^*GV)W \end{aligned}$$

It is straightforward that τ is a order serving map, thus by above lemma, it has a fixed point P . Let $V_1 = PV$ and $W_1 = (E - V^*PV)W$, then

$$\begin{aligned} V_1V_1^* &= P, & V_1^*V_1 &= V^*PV \\ W_1W_1^* &= F - P, & W_1^*W_1 &= E - V^*PV \end{aligned}$$

Thus $P \sim V^*PV$ and $F - P \sim E - V^*PV$, and by the additivity of equivalence, $F \sim E$. \square

6.3.2 Classification of Projections

Definition 1. A projection E in a von Neumann algebra \mathfrak{A} is finite if the only projection F in \mathfrak{A} s.t. $F \leq E$ and $F \sim F$ is $F = E$. Otherwise, E is infinite. A von Neumann algebra \mathfrak{A} is called finite or infinite if the identity is finite or infinite.

Remark. This definition generalize the concept of finiteness. It is useful. In fact, any projection containing in an abelian von Neumann algebra \mathfrak{A} is finite. Finite von Neumann algebras have similar properties to abelian von Neumann algebras. Therefore under this classification of projections, the concept of abelian von Neuman algebras can be generalized.

Proposition 2. 1) *If E is finite and $F \lesssim E$, then F is finite.*

2) *If \mathfrak{A} is a finite von Neumann algebra and E is a projection in \mathfrak{A} , then \mathfrak{A}_E is finite.*

3) *If $\{E_i\}$ is a family of finite projections in \mathfrak{A} s.t. $\{C_{E_i}\}$ is pairwise orthogonal, then $\sum_i E_i$ is a finite projection.*

4) *If $\{Z_i\}$ is a family of finite central projections, then $Z = \bigwedge_i Z_i$ is a finite central projection.*

5) If $\rho: \mathfrak{A} \rightarrow \mathfrak{B}$ is a $*$ -isomorphism and E is a finite projection in \mathfrak{A} , then $\rho(E)$ is a finite projection in \mathfrak{B} .

Proof. For 1), assume that $F \leq E$. If $F_1 \leq F$ and $F_1 \sim F$, then by the additivity,

$$F_1 + (E - F) \sim E$$

By finiteness, $F_1 + (E - F) = E$. Therefore, $F_1 = F$, i.e. F is finite. Then it is sufficient to prove that $F \sim E$ implies that F is finite. Let W be the partial isometry s.t. $W^*W = F$ and $WW^* = E$. If $F_1 \leq F$ and $F_1 \sim F$, then $E_1 = WF_1$ satisfies

$$E_1 \sim F_1 \sim F \sim E$$

Since E is finite, $E_1 = E$. Therefore, $F = W^*E = W^*WE_1 = F_1$. Thus F is finite. 2) holds by 1).

For 3), put $E = \sum_i E_i$. Let $F \leq E$ be a projection in \mathfrak{A} s.t. $F \sim_W E$. For any i ,

$$E - E_i \leq \sum_{j \neq i} C_{E_j} \Rightarrow (E - E_i)C_{E_i} = 0$$

Therefore, $E_i = EC_{E_i}$. Then for $W_i = WC_{E_i}$, $FC_{E_i} \sim_{W_i} EC_{E_i} = E_i$. Since $FC_{E_i} \leq EC_{E_i} = E_i$, $FC_{E_i} = E_i$. Therefore,

$$E = \sum_i E_i = \sum_i FC_{E_i} \leq F \leq E$$

So $F = E$ and E is finite.

For 4), let $F \leq Z$ and $F \sim Z$, then $Z_i F \sim Z_i Z = Z_i$. Since Z_i is finite, $Z_i F = Z_i$. Therefore, $F = Z$ and Z is finite.

5) is true by definition. \square

Corollary 3. If $\{\mathfrak{A}_i\}$ is a collection of von Neumann algebras, then $\bigotimes \mathfrak{A}_i$ is a finite von Neumann algebra if and only if each \mathfrak{A}_i is finite.

There are some similar properties of finite von Neumann algebra as abelian von Neumann algebra.

Proposition 4. \mathfrak{A} is finite if and only if every left (or right) invertible element of \mathfrak{A} is invertible.

Proof. Assume that every left (or right) invertible element of \mathfrak{A} is invertible. If $E \sim_W 1$. Since $WW^* = 1$, $W^*W = 1 = E$. Thus, \mathfrak{A} is finite. Conversely, if \mathfrak{A} is finite and $A, B \in \mathfrak{A}$ s.t. $AB = 1$, then $\text{ran } A = 1$. But $\overline{\text{ran}(A^*)} \sim \overline{\text{ran}(A)}$. Therefore, $\overline{\text{ran}(A^*)} = 1$ by the finiteness. So, $\ker A = 0$. And combining the fact $\text{ran } A = 1$, i.e A is surjection, A is invertible. \square

Proposition 5. Let \mathfrak{A} be a finite von Neumann algebra.

- 1) If E, E_1, F and F_1 are projections in \mathfrak{A} s.t. $E_1 \leq E$, $F_1 \leq F$ and $E \sim F$, $E_1 \sim F_1$, then

$$E - E_1 \sim F - F_1$$

- 2) If E and F are equivalent projections in \mathfrak{A} , then there is a unitary U in \mathfrak{A} s.t. $U^*EU \sim F$.

Proof. For 1), by the Comparison Theorem, there is a central projection Z s.t.

$$\begin{aligned} (E - E_1)Z &\lesssim (F - F_1)Z \\ (F - F_1)Z^\perp &\lesssim (E - E_1)Z^\perp \end{aligned}$$

Choosing F_2 s.t. $F_2 \leq (F - F_1)Z$ and $F_2 \sim (E - E_1)Z$. Therefore, $F_2 + F_1Z \leq FZ$ and by the additivity

$$F_2 + F_1Z \sim (E - E_1)Z + EZ \sim EZ \sim FZ$$

Since FZ is finite, $F_2 + F_1Z = FZ$. Therefore, $F_2 = (F - F_1)Z$. That is, $(E - E_1)Z \sim (F - F_1)Z$. And similarly, $(E - E_1)Z^\perp \sim (F - F_1)Z^\perp$.

For 2), by 1), $E \sim_W F$ implies that $E^\perp \sim_V F^\perp$. Then let $U = W + V$. U is a unitary and $U^*EU \sim F$. \square

Definition 6. 1) E is called an abelian projection in a von Neumann algebra if $E\mathfrak{A}E \cong \mathfrak{A}_E$ is an abelian von Neumann algebra.

- 2) An von Neumann algebra \mathfrak{A} is called discrete if for every nonzero central projection Z there is a nonzero abelian projection with $F \leq E$.
- 3) An von Neumann algebra \mathfrak{A} is continuous if it contains no nonzero abelian projections.
- 4) A projection E in \mathfrak{A} is called discrete or continuous if \mathfrak{A}_E is discrete or continuous.

Proposition 7. Let \mathfrak{A} be a von Neumann algebra.

- 1) A projection E in \mathfrak{A} is discrete if and only if for every central projection Z in \mathfrak{A} with $ZE \neq 0$, there is an abelian projection F in \mathfrak{A} with $F \leq ZE$.
- 2) A projection E in \mathfrak{A} is continuous if and only if there is no nonzero abelian projection F in \mathfrak{A} satisfying $F \leq E$.
- 3) Let \mathfrak{B} be a von Neumann algebra and $\rho: \mathfrak{A} \rightarrow \mathfrak{B}$ be a normal, surjective $*$ -homomorphism, and E be a projection in \mathfrak{A} . If E is continous (or discrete or abelian), then $\rho(E)$ is continous (or discrete or abelian).

- 4) If E is continous (or discrete or abelian) in \mathfrak{A} and F is a projection in \mathfrak{A} with $F \leq E$, then F is continous (or discrete or abelian).
- 5) If $\{E_i\}$ is a family of continous (or discrete or abelian) projections in \mathfrak{A} s.t. $\{C_{E_i}\}$ is pairwise orthogonal, then $\sum_i E_i$ is a continous (or discrete or abelian) projection.

Proof. For 1), it holds because that the center of \mathfrak{A}_E is \mathcal{Z}_E . Then for any $Z \in \mathcal{Z}$, there is an abelian projection $F \leq E$ in \mathfrak{A} with $FE \neq 0$ s.t. $EFE \leq EZE$, i.e. $F \leq ZE$. The converse is similar.

The proof of 2) is similar as 1).

For 3), it can assume ρ is a $*$ -isomorphism, then it follows by the definition.

If $E \sim F$, then there is a normal, surjective $*$ -homomorphism from \mathfrak{A}_E to \mathfrak{A}_F . Therefore, by 3), it sufficiently assumes $F \leq E$. Then since $F\mathfrak{A}F \subset E\mathfrak{A}E$, if E is abelian or continuous, F is clearly abelian or continuous. Assume E is discrete and $F \leq E$. Replacing \mathfrak{A}_E by \mathfrak{A} , i.e. assuming $E = 1$ and \mathfrak{A} is discrete. Let Z be a central projection with $ZF \neq 0$. Thus

$$ZC_F = C_{ZF} \neq 0$$

Then by the discreteness of \mathfrak{A} , there is an abelian projection $G \leq ZC_F$. Therefore,

$$C_G C_{ZF} = C_{ZF} \neq 0$$

Then there are $G_1 \leq G$ and $F_1 \leq ZF$ s.t. $G_1 \sim F_1$. Thus $F_1 \lesssim G$, i.e. F_1 is abelian. Then by 1), F is discrete.

Let E_i be abelian. Since $\{C_{E_i}\}$ is pairwise orthogonal, if $A, B \in \mathfrak{A}$, by the fact, $EAE = \sum_i (EAE)C_{E_i}$ because $C_E = \sum_i C_{E_i}$,

$$\begin{aligned} (EAE)(EAE) &= \left(\sum_i E_i A E_i \right) \left(\sum_j E_j B E_j \right) \\ &= \sum_i (E_i A E_i)(E_i B E_i) \\ &= \sum_i (E_i B E_i)(E_i A E_i) \\ &= (EBE)(EAE) \end{aligned}$$

Therefore, E is abelian.

Assume E_i is discrete. If Z is a central projection with $ZE \neq 0$, then for some i , $ZE_i \neq 0$. Then let $G_i \leq ZE_i$ be the abelian projection. If $ZE_i = 0$, let $G_i = 0$. Put $G = \sum_i G_i$ and since $C_{G_i} \leq C_{E_i}$, G is abelian. Therefore, E is discrete. Similarly, using 2), if E_i is continuous, E is continuous. \square

Remark. By 1) and 2), a projection in \mathfrak{A} is either discrete and continuous.

Corollary 8. *If $\{\mathfrak{A}_i\}$ is a collection of von Neumann algebras and $\mathfrak{A} = \bigoplus_i \mathfrak{A}_i$, then \mathfrak{A} is continuous (or discrete or abelian) if and only if each \mathfrak{A}_i is continuous (or discrete or abelian).*

If two projections E and F satisfy $E \lesssim F$, then $C_F \leq C_E$. But the converse is not always true. However, if E and F are abelian, the converse holds.

Lemma 9. *If E is an abelian projection and $F \leq E$, then $F = C_F E$.*

Proof. Since E is abelian, $E\mathfrak{A}E$ is abelian. In fact, $F\mathfrak{A}(E - F) = \{0\}$. It is because that for any $A \in \mathfrak{A}$,

$$FA(E - F) = F(EAE)(E - F) = F(E - F)(EAE) = 0$$

Therefore,

$$C_F = \overline{\mathfrak{A}F} \perp C_{E-F} = \overline{\mathfrak{A}(E - F)}$$

In particular, $C_F(E - F) = 0$. □

Proposition 10. *If E and F are abelian projections, then $F \lesssim E$ if and only if $C_F \leq C_E$.*

Proof. Assume $C_F \leq C_E$. By Comparison Theorem, there is a central projection Z s.t. $ZE \lesssim ZF$ and $Z^\perp F \lesssim Z^\perp E$. It can assume that $ZE \leq ZF$ by above proposition. Since ZF is abelian and by above lemma,

$$ZE = ZFC_{ZE} = ZFC_E = ZF$$

Therefore,

$$F = ZF + Z^\perp F = ZE + Z^\perp F \lesssim E \quad \square$$

6.3.3 Properties of Continuous Projections

For continuous projections, there is some decomposition like Lebesgue Decomposition for general measures.

Lemma 1. *If E is not an abelian projection, then there are projections E_1, E_2 in \mathfrak{A} that are dominated by E and s.t. $E_1 \perp E_2$ and $E_1 \sim E_2$.*

Proof. Since E is not abelian, there is a projection $F \leq E$ that does not belong to the center of $E\mathfrak{A}E$, $\mathcal{Z}E$.

Claim: $F\mathfrak{A}(E - F) \neq \{0\}$.

If for any $A \in \text{Re } \mathfrak{A}$, $FA(E - F) = 0$. Then $FAE = FAF$ and $EAF = FAF$. Therefore,

$$F(EAE) = FAE = FAF = EAF = (EAE)F$$

That is F is in the center of $E\mathfrak{A}E$, which is a contradiction.

$F\mathfrak{A}(E - F) \neq \{0\}$ implies that there are $E_1 \leq F$ and $E_2 \leq E - F$ s.t. $E_1 \sim E_2$. □

Theorem 2. *If E is continuous, then there are projections E_1 and E_2 in \mathfrak{A} s.t. $E_1 \perp E_2$, $E_1 \sim E_2$ and $E = E_1 + E_2$.*

Proof. Consider the collection of all sets of pairs $\{E_{1i}, E_{2i}\}$ s.t. for all i ,

$$E_{1i} \sim E_{2i}, E_{1i}, E_{2i} \leq E_{ni} \perp E_{mj}$$

for $i \neq j$ and $m, n = 1, 2$. Then by above lemma this set is not empty and then by Zorn's Lemma, there is a maximal element $\{E_{1i}, E_{2i}\}$ and put

$$E_1 = \sum_i E_{1i}, E_2 = \sum_i E_{2i}$$

By the additivity, $E_1 \sim E_2$ and clearly, $E_1 \perp E_2$. And by the maximality and the continuity of E , $E = E_1 + E_2$. \square

Corollary 3. *The projection E is continuous if and only if for every $F \leq E$, there are projections F_1 and F_2 in \mathfrak{A} s.t. $F_1 \perp F_2$, $F_1 \sim F_2$ and $F = F_1 + F_2$.*

Proof. If E is continuous, so is each its subprojection, thus by above theorem it is true. And the converse is trivial. \square

6.4 Type I Von Neumann Algebras

6.4.1 Classification of Von Neumann Algebras

Then by using the classifications of projections, von Neumann algebras can be classified as following categories.

Definition 1. A von Neumann algebra \mathfrak{A} is:

- 1) Type I_n if it is finite and discrete.
- 2) Type I_∞ if it is infinite and discrete.
- 3) Type II_1 if it is finite and continuous.
- 4) Type II_∞ if it is infinite and continuous and every nonzero central projection dominates a nonzero finite projections.
- 5) Type III if it contains no nonzero finite projections.

Proposition 2. *Let \mathfrak{A} be a von Neumann algebra and $E \in \mathfrak{A}$ be a projection.*

- 1) *If \mathfrak{A} is Type I (II, or III), then \mathfrak{A}_E is Type I (II, or III).*
- 2) *If \mathfrak{A} is Type I_n , then \mathfrak{A}_E is Type I_n .*
- 3) *If \mathfrak{A} is Type II_1 , then \mathfrak{A}_E is Type II_1 .*

Proof. Just need to prove the cases of Type II and Type III. Let Z be a central projection in \mathfrak{A} s.t. $EZ \neq 0$. If \mathfrak{A} is Type II, then there is a nonzero finite projection F with $F \leq Z$. Thus

$$EZ F = EZ \neq 0$$

Then there are $E_1 \leq EZ$ and $F_1 \leq F$ s.t. $E_1 \sim F_1$. Therefore, E_1 is finite and thus \mathfrak{A}_E is Type II. If \mathfrak{A} is Type III, then \mathfrak{A} contains no nonzero finite projection, thus \mathfrak{A}_E contains no nonzero finite projection, and continuity of \mathfrak{A}_E is clear. Therefore, \mathfrak{A}_E is Type III. \square

Proposition 3. *Let $\{A_i\}$ be a class of von Neumann algebras and put $\mathfrak{A} = \bigoplus_i \mathfrak{A}_i$.*

- 1) \mathfrak{A} is Type I_n if and only if each \mathfrak{A}_i is Type I_n .
- 2) \mathfrak{A} is Type I_∞ if and only if each \mathfrak{A}_i is Type I and at least one of them is Type I_∞ .
- 3) \mathfrak{A} is Type II_1 if and only if each \mathfrak{A}_i is Type II_1 .
- 4) \mathfrak{A} is Type II_∞ if and only if each \mathfrak{A}_i is Type II and at least one of them is Type II_∞ .
- 5) \mathfrak{A} is Type III if and only if each \mathfrak{A}_i is Type III.

Theorem 4. *If \mathfrak{A} is a von Neumann algebra, then there are unique pairwise orthogonal central projections Z_1, \dots, Z_5 with $\sum_{i=1}^5 Z_i = 1$ s.t.*

- 1) $\mathfrak{A}Z_1$ is Type I_n .
- 2) $\mathfrak{A}Z_2$ is Type I_∞ .
- 3) $\mathfrak{A}Z_3$ is Type II_1 .
- 4) $\mathfrak{A}Z_4$ is Type II_∞ .
- 5) $\mathfrak{A}Z_5$ is Type III.

Proof. Define the central projections as

$$\begin{aligned} H &= \sup \{Z : Z \text{ is a central projection, } \mathfrak{A}Z \text{ is Type I}\} \\ K &= \sup \{Z : Z \text{ is a central projection, } \mathfrak{A}Z \text{ is Type II}\} \end{aligned}$$

Then it is easy to see $\mathfrak{A}H$ is Type I and $\mathfrak{A}K$ is Type II and $HK = 0$. Put

$$Z_5 = 1 - (H + K)$$

Check: $\mathfrak{A}Z_5$ is continuous.

Assume there is an abelian projection P in $\mathfrak{A}Z_5$. Clearly, its central cover

$C_P \leq Z_5$. If G is any nonzero central projection in $\mathfrak{A}C_P$. Then G is also a central projection in \mathfrak{A} . Moreover, $GP \neq 0$, and $(GP)\mathfrak{A}(GP) = G(P\mathfrak{A}P)$ is abelian. That means $\mathfrak{A}C_P$ is Type I. Thus $C_P \leq H$. Therefore, $C_P = 0$. Then $P = 0$. Since P is arbitrary, $\mathfrak{A}Z_5$ is continuous.

Check: $\mathfrak{A}Z_5$ is Type III.

Suppose there is a nonzero finite projection P in $\mathfrak{A}Z_5$. Similarly, $\mathfrak{A}C_P$ is continuous. If G is any central projection of $\mathfrak{A}C_P$, then $G \in \mathcal{Z}$ and $G \leq C_P$. Thus $GC_P \neq 0$ and GP is finite. Therefore, $\mathfrak{A}C_P$ is Type II. Thus $C_P \leq K$, then $C_P = 0$ i.e. $P = 0$.

Then define $F = \sup \{Z : Z \text{ is a finite central projection}\}$. Put $Z_1 = HF$, $Z_2 = HF^\perp$, $Z_3 = KF$ and $Z_4 = KF^\perp$. \square

Corollary 5. *If \mathfrak{A} is a factor, then \mathfrak{A} is exactly one of the Type I_n , I_∞ , II_1 , II_∞ and III.*

6.4.2 The Structure of Type I Algebras

Definition 1. A projection is faithful if its central cover is identity.