

Multiple recurrence and its Noncommutative Extensions

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DECLARATION

I hereby declare that I am the sole author of this thesis in partial fulfillment of the requirements for a Integrated MSc. degree from National Institute of Science Education and Research (NISER). I authorize NISER to lend this thesis to other institutions or individuals for the purpose of scholarly research.

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Date:

The thesis work reported in the thesis entitled “**Multiple recurrence and its Noncommutative Extensions**” was carried out under my supervision, in the School of Mathematical Sciences at NISER, Bhubaneswar, India.

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ABSTRACT

In the first half, we present the proof of the multiple recurrence theorem. We begin with the basics of ergodic theory. We discuss ergodicity, weak-mixing systems and the ergodic theorems: the mean ergodic theorem and pointwise ergodic theorem. We then delve into multiple recurrence. And we give the proof of Szemerédi's theorem as an application.

In the second half, we explore the extension of recurrence theorems from measure preserving systems to their noncommutative setting, i.e., C^* -dynamical systems. We first set up the theory of C^* -algebras and von Neumann algebras, with some emphasis on W^* -algebra. Then, we discuss Poincaré recurrence for C^* -dynamical systems and von Neumann algebras. We conclude with the proof of multiple recurrence theorem for almost-periodic C^* -dynamical systems.

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Chapter 1

Preliminaries

1.1 Measure Theory

1.1.1 Basics

We begin with some basic definitions and results from measure theory.

Definition 1.1.1.

Algebra: on set X is a collection of subsets $\mathcal{A} \subset \mathcal{P}(X)$ such that $X \in \mathcal{A}$ and is closed under complementation and finite union.

σ -algebra: on set X is a collection of subsets $\mathcal{B} \subset \mathcal{P}(X)$ such that $X \in \mathcal{B}$ and is closed under complementation and countable union.

σ -algebra: generated by $S \subset X$, $\sigma(S)$ is the smallest σ -algebra containing S .

Monotone class: is a collection of subsets $\mathcal{M} \subset \mathcal{P}(X)$ that is closed under union of increasing sequences and intersection of decreasing sequences.

Measurable space: is a pair (X, \mathcal{B}) , where X is a set and \mathcal{B} is a σ -algebra on X .

Measure: on (X, \mathcal{B}) is a *countably additive* function $\mu : X \rightarrow [0, \infty)$: if $\{E_i\}$ is a countable disjoint collection in \mathcal{B} , then $\mu(\bigcup E_i) = \sum \mu(E_i)$.

Probability space: is a triple (X, \mathcal{B}, μ) , where X, \mathcal{B} is measurable space and μ is a measure on it such that $\mu(X) = 1$ (μ is called a probability measure).

Measure-preserving transformation: is a measurable function $T : X \rightarrow X$: if T such that $\mu \circ T^{-1} = \mu$.

Measure-preserving system: (MPS) is a quadruple (X, \mathcal{B}, μ, T) , where (X, \mathcal{B}, μ) is a probability space and $T : X \rightarrow X$ is a measure-preserving transformation.

Theorem 1.1.2 (Monotone class theorem). *If \mathcal{A} is an algebra, then $M(\mathcal{A}) = \sigma(\mathcal{A})$, where $M(\mathcal{A})$ is the smallest monotone class containing \mathcal{A} .*

Definition 1.1.3 (Push-forward measure). Let (X, \mathcal{B}, μ) be a measure space, (Y, \mathcal{C}) and $\phi : X \rightarrow Y$ be a measurable map. Then the *push-forward measure* is a measure on Y defined by

$$\mu \circ \phi^{-1}(E) = \mu(\phi^{-1}(E)).$$

Theorem 1.1.4 (Change of variables). *Let $\phi : (X, \mathcal{B}, \mu) \rightarrow (Y, \mathcal{C})$ be a measurable map and f be a real-valued measurable function on Y , then*

$$\int_Y f d(\mu \circ \phi^{-1}) = \int_X f \circ \phi d\mu$$

1.1.2 Conditional Expectation

Definition 1.1.5 (Absolutely continuous). Let ν, μ be measures on (X, \mathcal{B}) . ν is absolutely continuous with respect to μ if

$$\mu(A) = 0 \implies \nu(A) = 0 \quad \forall A \in \mathcal{B}.$$

Theorem 1.1.6 (Radon-Nikodym theorem). Let ν, μ be σ -finite measures on (X, \mathcal{B}) such that $\nu \ll \mu$. There exists a measurable function $f : X \rightarrow [0, \infty)$ such that

$$\nu(A) = \int_A f d\mu \quad \forall A \in \mathcal{B}.$$

Definition 1.1.7 (Measure-preserving system). A quadruple (X, \mathcal{B}, μ, T) , where (X, \mathcal{B}, μ) is a probability space and $T : X \rightarrow X$ is a *measure-preserving transformation*.

Definition 1.1.8 (Conditional expectation). Let (X, \mathcal{B}, μ) be a probability space, $f \in L^1(X, \mathcal{B}, \mu)$, and $\mathcal{A} \subset \mathcal{B}$ be a sub- σ -algebra. The *conditional expectation of f given \mathcal{A}* , denoted $\mathbb{E}(f|\mathcal{A})$, is a \mathcal{A} -measurable function such that

$$\int_A \mathbb{E}(f|\mathcal{A}) d\mu = \int_A f d\mu \quad \forall A \in \mathcal{A}.$$

Such a function exists and is a.e. unique due to the Radon-Nikodym theorem. We first consider non-negative f . Define

$$\nu(A) = \int_A f d\mu \quad \text{for } A \in \mathcal{A}.$$

Then $\nu(A)$ is a finite measure on \mathcal{A} and is absolutely continuous with respect to $\mu|_{\mathcal{A}}$. By Radon-Nikodym theorem, there is a function $g \in L^1(X, \mathcal{A}, \mu)$ such that

$$\nu(A) = \int_A g d\mu \quad \forall A \in \mathcal{A}.$$

We denote $g = \mathbb{E}(f|\mathcal{A})$.

For general $f \in L^1$, we write $f = f^+ - f^-$ and get $\mathbb{E}(f|\mathcal{A}) = \mathbb{E}(f^+|\mathcal{A}) - \mathbb{E}(f^-|\mathcal{A})$.

Example. $\{A_i\}_{i=0}^n$ is a finite partition of X , and \mathcal{A} is the σ -algebra generated by it. Since $\mathbb{E}(f|\mathcal{A})$ is \mathcal{A} -measurable, it must be a linear combination of 1_{A_i} 's,

$$\mathbb{E}(f|\mathcal{A}) = \sum_{i=0}^n a_i 1_{A_i}.$$

We can compute,

$$a_i = \frac{\int_{A_i} f d\mu}{\mu(A_i)}.$$

If $f = 1_A$ for $A \subset X$, then $a_i = \frac{\mu(A_i \cap A)}{\mu(A_i)}$. This is equivalent to the definition of ‘conditional probability of A given B ’ for events A, B from elementary probability theory.

Theorem 1.1.9. *Conditional expectation has the following properties:*

- (i) $\mathbb{E}(af + bg|\mathcal{A}) = a\mathbb{E}(f|\mathcal{A}) + b\mathbb{E}(g|\mathcal{A})$ a.e., where $a, b \in \mathbb{R}$.
- (ii) $f \leq g$ a.e., then $\mathbb{E}(f|\mathcal{A}) \leq \mathbb{E}(g|\mathcal{A})$ a.e.
- (iii) $|\mathbb{E}(f|\mathcal{A})| \leq \mathbb{E}(|f||\mathcal{A})$

Proof. (i) For any $A \in \mathcal{A}$,

$$\int_A a\mathbb{E}(f|\mathcal{A}) + b\mathbb{E}(g|\mathcal{A}) d\mu = \int_A af + bg d\mu = \int_A \mathbb{E}(af + bg|\mathcal{A}).$$

(ii) For any $A \in \mathcal{A}$,

$$\int_A \mathbb{E}(f|\mathcal{A}) - \mathbb{E}(g|\mathcal{A}) d\mu = \int_A f - g d\mu \geq 0.$$

Thus, $\mathbb{E}(f|\mathcal{A}) - \mathbb{E}(g|\mathcal{A}) \geq 0$ a.e.

(iii) Follows from (ii) by the fact $f \leq |f|$. □

1.2 Functional Analysis

We record some important results from functional analysis.

Theorem 1.2.1 (Banach-Alaoglu theorem). *X is a normed space. Then the unit ball of X^* is a compact Hausdorff space.*

Theorem 1.2.2 (Hahn-Banach extension theorem). *X is a normed space, and $X_0 \subset X$ be a subspace. If $\phi : X_0 \rightarrow \mathbb{C}$ is a bounded linear map, then there exists a bounded linear map $\tilde{\phi} : X \rightarrow \mathbb{C}$ such that,*

$$(i) \quad \tilde{\phi}|_{X_0} = \phi,$$

$$(ii) \quad \|\tilde{\phi}\| = \|\phi\|.$$

Theorem 1.2.3 (Hahn-Banach separation theorem). *Let X be a normed linear space, and $A, B \subset X$ be nonempty disjoint convex subsets.*

(i) *If A is open, then there exists a continuous linear functional $f : X \rightarrow \mathbb{R}$, $\alpha \in \mathbb{R}$ such that*

$$f(a) \leq \alpha \leq f(b), \quad \forall a \in A, \quad \forall b \in B.$$

(ii) *If A is compact and B is closed, then there exists a continuous linear functional $f : X \rightarrow \mathbb{R}$, $\alpha \in \mathbb{R}$ such that*

$$f(a) < \alpha < f(b), \quad \forall a \in A, \quad \forall b \in B.$$

Theorem 1.2.4 (Open mapping theorem). *Let $T : X \rightarrow Y$ be a surjective bounded linear map between two Banach spaces, X and Y , then T is an open map.*

Theorem 1.2.5 (Closed graph theorem). *Let $T : X \rightarrow Y$ be a linear map between two Banach spaces, X and Y , then T is continuous if and only if the graph $G(T) = \{(x, Tx) \mid x \in X\}$ is closed in $X \times Y$.*

Theorem 1.2.6 (Uniform boundedness theorem). *Let X be a Banach space, Y be a normed space and $\{T_j\}$ be a family of bounded operators from X to Y , then the family $\{T_j\}$ is pointwise bounded if and only if the family is uniformly bounded, i.e., $\sup_j \|T_j\| < \infty$.*

Theorem 1.2.7 (Spectral theorem for compact operators). *Let $x \in \mathcal{B}(\mathcal{H})$ be a compact self-adjoint operator. Then there exists $\lambda_1, \lambda_2, \dots \in \mathbb{R}$ such that*

$$1. \quad |\lambda_1| \leq |\lambda_2| \leq \dots \text{ and } \lim_i \lambda_i = 0,$$

$$2. \quad x = \sum_{i=1}^{\infty} \lambda_i p_i, \text{ where } p_i \text{ is the orthogonal projection onto } \ker(x - \lambda_i).$$

Chapter 2

Ergodic Theory

2.1 Recurrence

Recurrence is a basic and deep property of all measure-preserving system.

Lemma 2.1.1. *If $\{E_i\}$ is a countable collection in \mathcal{B} such that $\mu(A_i \cap A_j) = 0$ when $i \neq j$, then $\mu(\bigcup E_i) = \sum \mu(E_i)$.*

Proof. Let $N = \bigcup_{i \neq j} (A_i \cap A_j)$, then $\mu(N) = 0$. Let $B_i = A_i \setminus N$, then

$$\begin{aligned}\mu(\bigcup B_i) &= \mu(\bigcup (A_i \setminus N)) = \mu((\bigcup A_i) \setminus N) = \mu(\bigcup A_i), \\ \mu(\bigcup B_i) &= \sum \mu(B_i) = \sum \mu(A_i).\end{aligned}$$

□

Lemma 2.1.2. *Let A be a measurable set with positive measure, then there is an n such that $\mu(A \cap T^{-n}A) > 0$.*

Proof. Take $k > 1/\mu(A)$, and consider the sets $A, T^{-1}A, \dots, T^{-k}A$. $T^{-i}A$ are not almost pairwise disjoint, since if they were, then by using lemma 2.1.1,

$$1 = \mu(X) \geq \mu(\bigcup_{i=1}^k T^{-i}A) = \sum_{i=1}^k \mu(A) > 1.$$

So, there are $1 \leq i < j \leq k$ such that $\mu(T^{-i}A \cap T^{-j}A) > 0$. Then $\mu(T^{-i}(A \cap T^{-(j-i)}A)) = \mu(A \cap T^{-(j-i)}A) > 0$. □

Theorem 2.1.3 (Poincare recurrence theorem). *Let A be a measurable set with positive measure, then almost every point of A returns to A , i.e., there is a set $E \subset A$, $\mu(E) = 0$ such that if $x \in A \setminus E$, then there exist $n \in \mathbb{N}$, $T^n x \in A$. Furthermore, the points return infinitely often, i.e., there are infinitely many n such that $T^n x \in A$.*

Proof. Consider the set of points of A that do not return to A , $E = A \setminus \bigcup_{n \in \mathbb{N}} T^{-n}A$. Then, $E \cap T^{-n}E \subset A \cap T^{-n}E = \emptyset$ for all $n \in \mathbb{N}$. So, by the previous lemma, $\mu(E) = 0$, and every point of E returns to A at least once.

To show they return infinitely often, consider the set of points of A that return only finitely many points. Note, if $x \in A$ returns to A k times then there are $n_1 < \dots < n_k$ such that $T^{n_k} x \in E$. So, $\bigcup_{n_1 < \dots < n_k} T^{-n_k}E$, where the union is over all k -tuples, $\{n_1 < \dots < n_k\}$ consists of the points of A that return exactly k many times. Thus, the set $\bigcup_{k \in \mathbb{N}} \bigcup_{n_1 < \dots < n_k} T^{-n_k}E$ are the set of points of A that return finitely many times, and has measure 0 since $\mu(E) = 0$. □

2.2 Ergodicity

2.2.1 Ergodicity

Definition 2.2.1 (Invariant set). A measurable set A is an *invariant set* if $T^{-1}A = A$.

Definition 2.2.2 (Invariant function). A measurable function f is an *invariant function* if $f \circ T = f$.

Lemma 2.2.3. f is bounded, \mathcal{B} -measurable and A is a T -invariant set, then $\int_A f \circ T d\mu = \int_A f d\mu$.

Proof.

$$\int_A f \circ T d\mu = \int_X (f \circ T) 1_{T^{-1}A} d\mu = \int_X (f 1_A) \circ T d\mu = \int_X f 1_A d(\mu \circ T^{-1}) = \int_A f d\mu.$$

□

Note, if A is an invariant set, then $X \setminus A$ is also an invariant set.

Definition 2.2.4 (Ergodic system). A measurable-preserving system (X, \mathcal{B}, μ, T) is *ergodic* if there are no non-trivial invariant sets, i.e., if A is an invariant set, then $\mu(A) = 0$ or 1 .

Ergodicity is an irreducibility condition: if T is non-ergodic, then it can be reduced to simpler transformations $T|_A$ and $T|_{X \setminus A}$, since A and $X \setminus A$ do not interact, i.e., orbits in one of them do not enter the other.

There are several ways of defining ergodicity.

Theorem 2.2.5. (X, \mathcal{B}, μ, T) is a MPS, then the following are equivalent,

- (i) T is ergodic.
- (ii) if $T^{-1}A = A \mod\text{-}\mu$ [i.e. $\mu(T^{-1}A \triangle A) = 0$], then $\mu(A) = 0$ or 1 .
- (iii) if $f : X \rightarrow \mathbb{R}$ is T -invariant, then f is constant a.e.
- (iv) if $f \in L^1$ such that $f \circ T = f$ a.e., then f is constant a.e.

Proof.

- (i) \iff (iii): If T is not ergodic, there is a non-trivial invariant set A . Then 1_A is a measurable function such that $1_A \circ T = 1_{T^{-1}A} = 1_A$ which is not constant a.e.

Now, assume T is ergodic and $f \circ T = f$ a.e. For $k \in \mathbb{Z}$ and $n \in \mathbb{N}$, define

$$X(k, n) = \{x \in X \mid \frac{k}{2^n} \leq f(x) \leq \frac{k+1}{2^n}\}.$$

Since, $T^{-1}X(k, n) \triangle X(k, n) \subset \{x \mid f(x) \neq f(Tx)\}$, it must have measure 0 or 1. For a fixed n , $X(k, n)$ are disjoint and its union covers X , $\mu(X(k_n, n)) = 1$ for exactly one k_n . Now consider the set $\bigcap_n X(k_n, n)$. It has measure 1 and we can show that f is constant on it.

- (ii) \iff (iv): Same as previous.

- (iii) \iff (iv): If f is a T -invariant function, we construct integrable functions $f\chi_n$, where χ_n is the characteristic function on $\{|f| < n\}$. Then $f\chi_n$ are constant a.e., and since they have non-trivial intersection of domains, for all n the constants are same. As $f\chi_n \rightarrow f$ pointwise, we get that f is constant a.e.

If $f \in L^1$ such that $f \circ T = f$ a.e., let $g(x) = \limsup f(T^n(x))$, then g is measurable and $g \circ T = g$. So g is constant a.e. We are done if we show $f = g$ a.e.

Note, $f(T^{n+1}x) = f(T^n x) \forall n \implies f(T^n x) = f(x) \forall n \implies g(x) = f(x)$,
 and $f(T^{n+1}x) = f(T^n x) \iff T^n x \in \{fT = f\} \iff x \in T^{-n}\{fT = f\}$.
 So, $\{g = f\} \supset \bigcap_n T^{-n}\{fT = f\}$, which being the intersection of sets with measure 1, has measure 1. □

Theorem 2.2.6. (X, \mathcal{B}, μ, T) is a MPS, then the following are equivalent,

- (i) T is ergodic.
- (ii) if $A \in \mathcal{B}, \mu(A) > 0$, then $\bigcup_{n=m}^{\infty} T^{-n}A = X \text{ mod-}\mu \forall m$.
- (iii) if $A, B \in \mathcal{B}, \mu(A), \mu(B) > 0$, then $\mu(T^{-n}A \cap B) > 0$ for infinitely many n .

Proof.

- (i) \implies (ii): Let $A' = \bigcup_{n=m}^{\infty} T^{-n}A$, then $T^{-1}A' = \bigcup_{n=m+1}^{\infty} T^{-n}A \subset A'$. And since $\mu(T^{-1}A) = \mu(A)$, we have $T^{-1}A' = A' \text{ mod-}\mu$. By ergodicity, $A' = X \text{ mod-}\mu$.
- (ii) \implies (iii): We have $\mu(\bigcup_{n=1}^{\infty} T^{-n}A) = 1$. So, $0 < \mu(B) = \mu(B \cap \bigcup_n T^{-n}A) = \mu(\bigcup_n (B \cap T^{-n}A))$. Hence, for some n_0 , $\mu(B \cap T^{-n_0}A) > 0$. We can find another such n by taking the union from n_0 instead of 1 in the first line.
- (iii) \implies (i): Assume T is not ergodic and let A be a non-trivial invariant subset, then $A \cap (X \setminus A) = \emptyset \implies (T^{-n}A) \cap X \setminus A = \emptyset \forall n \implies \mu(T^{-n}A \cap (X \setminus A)) = 0 \forall n$, which is a contradiction. □

Example. On S^1 with the Lebesgue measure, the rotation $T_\alpha(z) = ze^{2\pi i\alpha}$ is ergodic if and only if α is irrational.

Proof. If $\alpha = p/q \in \mathbb{Q}$, $f(z) = z^q$ is a such that it is non-constant a.e, T_α -invariant function. Hence, T_α is not ergodic.

If $\alpha \in \mathbb{Q}^c$: Let f be a T_α -invariant function in L^1 . Consider the Fourier series, $f = \sum a_n \chi_n$, where $\chi_n(z) = z^n$ are the characters on S^1 . Then,

$$f \circ T_\alpha(z) = \sum a_n e^{2\pi i n \alpha} z^n = \sum a_n z^n = f(z).$$

By uniqueness of Fourier coefficients, for $n \neq 0$, $a^n e^{2\pi i n \alpha} = a_n$. Which implies that $a_n = 0$, when $n \neq 0$, and that f is not constant a.e. □

Let (X, \mathcal{B}, μ, T) be an ergodic measure-preserving system and let $A \in \mathcal{B}$ with positive measure. By ergodicity, $\bigcup_{n=1}^{\infty} T^{-n}A = X \text{ mod-}\mu$, which says that a.e. $x \in X$ returns to A . Kac's formula gives us the average return time.

Define the function $r_A(x) : X \rightarrow \mathbb{N}$ by $r_A(x) = \min\{r \geq 1 \mid T^r x \in A\}$. Since $n_A((-\infty, k]) = \bigcup_{i=0}^k T^{-i}A$ is measurable for every $k \in \mathbb{N}$, r_A is measurable.

Theorem 2.2.7 (Kac's return time formula). *Let (X, \mathcal{B}, μ, T) be a MPS such that T is invertible. Then*

$$\int_A r_A d\mu = 1.$$

Proof. Let $A_n = A \cap \{r_A = n\}$, then on A , $r_A = \sum_{n=1}^{\infty} n 1_{A_n}$. So,

$$\int_A r_A d\mu = \sum_{n=1}^{\infty} n \mu(A_n) = \sum_{n=1}^{\infty} \sum_{m=1}^n \mu(T^m A_n).$$

We show $\{T^m A_n \mid n \in \mathbb{N}, m = 1, \dots, n\}$ are disjoint and their union has full measure. Then, the result follows from countable additivity of μ .

To show: $\mu(\bigcup_{n=1}^{\infty} \bigcup_{m=1}^n T^m A_n) = 1$.

Proof. By ergodicity, a.e. $x \in X$, there is a smallest $m \geq 1$ such that $y := T^{-m}x \in A$. Let $n = r_A(y)$. Then, if $m > n$, $T^n y = T^{-(m-n)}x \in A$, but $m - n < m$, which is a contradiction of the minimality. Thus, $m \leq n$, and $x \in T^m A_n$. \square

To show: $\{T^m A_n \mid n \in \mathbb{N}, m = 1, \dots, n\}$ are disjoint.

Proof. We show if $x \in T^{m'} A_{n'}$, then $(m', n') = (m, n)$. First, $T^{-m'} x \in A_{n'} \subset A$. So, by minimality, $m \leq m'$. If $m < m'$, since $T^{-m} x = T^{m'-m}(T^{-m'} x) \in A_{n'} \subset A$, $n' = r_A(T^{-m'} x) \leq m' - m < m'$, which is a contradiction. Thus, $m = m'$. Now, if $x \in T^m A_n \cap T^{m'} A_{n'}$, observe if $n \neq n'$, $A_n \cap A_{n'} = \emptyset \implies T^m A_n \cap T^{m'} A_{n'} = \emptyset$. Hence, $n = n'$. \square

\square

2.2.2 Ergodic decomposition

We have seen that ergodic systems are non-decomposable. In this section, we look at how non-ergodic systems can be decomposed into ergodic ones. To do so, we will use measure disintegration.

Definition 2.2.8 (Probability kernel). Let (X, \mathcal{B}) and (Y, \mathcal{C}) be two measurable spaces and $\{\theta_x\}_{x \in X}$ be a family of probability measures on Y , then $\{\theta_x\}_{x \in X}$ is called a *probability kernel from (X, \mathcal{B}) to (Y, \mathcal{C})* if for each $E \in \mathcal{C}$, the map $x \mapsto \theta_x(E)$ is \mathcal{B} -measurable.

Example. Every measurable map $\phi : (X, \mathcal{B}) \rightarrow (Y, \mathcal{C})$ induces a kernel by $x \mapsto \delta_{\phi(x)}$, since for $E \in \mathcal{C}$, the map $x \mapsto \delta_{\phi(x)}(E) = 1_E(\phi(x))$ is a composition of measurable functions.

Definition 2.2.9 (Measure integration). Let (X, \mathcal{B}, μ) be a measure space, (Y, \mathcal{C}) be a measurable space and $\{\theta_x\}_{x \in X}$ be a kernel from (X, \mathcal{B}) to (Y, \mathcal{C}) , then we can define a probability measure on Y ,

$$\nu(E) = \int \nu_x(E) d\mu(x).$$

ν is a measure by an application of monotone convergence theorem: for positive measurable functions $\{f_n\}_{n=1}^{\infty}$, $\int \sum_{n=1}^{\infty} f_n = \sum_{n=1}^{\infty} \int f_n$.

We now look at the inverse problem and study how we can decompose a given measure as an integral of other measures. We can do this on some special spaces.

Definition 2.2.10 (Standard measurable space). A measurable space (X, \mathcal{B}) is *standard* if there exists a complete and separable metric on X for which \mathcal{B} is the Borel σ -algebra.

We will use the following result from advanced measure theory, that classifies standard measurable spaces, by cardinality.

Theorem 2.2.11. *There are three standard measurable spaces, up to measurable isomorphism: finite discrete space, countable discrete space and $[0, 1]$ with the usual Borel σ -algebra.*

Example. These standard measurable space are isomorphic: $[0, 1]$ with the Borel σ -algebra, $S^{\mathbb{Z}}$ (S is a finite set) with the product Borel σ -algebra.

Definition 2.2.12 (Measure disintegration). Let (X, \mathcal{B}, μ) be a measure space and $\mathcal{F} \subset \mathcal{B}$ a sub- σ -algebra, then a kernel from (X, \mathcal{F}) to (X, \mathcal{B}) , $\{\theta_x\}_{x \in X}$ is called the *disintegration of μ over \mathcal{F}* if

1. $\mu(E) = \int \theta_x(E) d\mu(x), \forall E \in \mathcal{B}$.

2. if $f : X \rightarrow \mathbb{C}$ is a bounded \mathcal{B} -measurable function, then

$$\mathbb{E}(f|\mathcal{F})(x) = \int f(t) d\theta_x(t) \text{ } \mu\text{-a.e.} \quad (2.1)$$

We begin by examining the simplest case when $\mathcal{F} = \sigma\{A_1, \dots, A_n\}$ is a finite σ -algebra, where $\{A_i\}_{i=1}^n$ are measurable sets that partition A . Denote $A(x)$ to be the unique A_i containing x , and we can define the kernel $\theta_x = \frac{\mu|_{A(x)}}{\mu(A(x))}$. Then,

$$\int \theta_x(E) d\mu(x) = \int \frac{\mu|_{A(x)}(E)}{\mu(A(x))} d\mu(x) = \sum_{i=1}^n \int \frac{\mu(E \cap A_i)}{\mu(A_i)} d\mu|_{A_i}(x) = \mu(E).$$

So, θ_x is a disintegration of μ over \mathcal{F} .

Writing $\theta_x(E) = \mathbb{E}(1_E|\mathcal{F})(x)$, we can extend this to when $\mu(E) = 0$. But since, conditional expectation is well-defined a.e., this works when there are countably many E . We deal with this technicality in the next theorem, by first finding a pre-measure on a countable algebra.

Theorem 2.2.13. *When (X, \mathcal{B}, μ) is a measure space and $\mathcal{F} \subset \mathcal{B}$ a sub- σ -algebra, there is a disintegration of μ over \mathcal{F} . Further, the kernel is unique: if θ'_x is another kernel, then $\theta_x = \theta'_x$ μ -a.e.*

Proof. We provide a proof for when $X = S^{\mathbb{Z}}$, $S = \{1, \dots, n\}$ is a finite set, with the Borel σ -algebra on the product topology.

Let $\pi_n : S^{\mathbb{Z}} \rightarrow S^{2n+1}$ be the projection into $-n^{\text{th}}$ to n^{th} coordinates, with the discrete topology on S^{2n+1} . Define $\mathcal{A}_n := \{\pi_n^{-1}(E) \mid E \subset A^{2n+1}\}$. This is an algebra and $\mathcal{A}_n \subset \mathcal{A}_{n+1}$. So, $\mathcal{A} := \bigcup_{n=1}^{\infty} \mathcal{A}_n$ is also an algebra. Now, we can define

$$\theta_x(E) = \mathbb{E}(1_E|\mathcal{F})(x), \text{ for } E \in \mathcal{A}.$$

Since \mathcal{A} is countable, it is well-defined on a set of full measure and $0 \leq \theta_x(E) \leq 1$, call it X_0 .

Claim: θ_x is a pre-measure on \mathcal{A} .

Proof. Clearly, $\theta_x(X) = 1$, and is finitely additive. Instead of showing countable additivity, we show that finite additivity is enough, since no element in \mathcal{A} can be written as a countable union. Let $A \in \mathcal{A}$ such that $A = \bigsqcup_{i=1}^{\infty} A_i$, $A_i \in \mathcal{A}$. It follows that infinitely many A_i are empty from the fact A is compact and A_i are disjoint. \square

Now, we can extend θ_x to a probability measure, which we also call θ_x , on $\sigma(\mathcal{A})$, by Caratheodory extension theorem. And $\sigma(\mathcal{A}) = \mathcal{B}$. Indeed, we can check that \mathcal{A} forms a basis on the product topology.

Claim: The map $X_0 \ni x \mapsto \theta_x(E)$ is \mathcal{F} -measurable and $\theta_x(E) = \mathbb{E}(1_E|\mathcal{F})(x)$, for $E \in \mathcal{B}$.

Proof. Let $\mathcal{C} = \{B \in \mathcal{B} \mid x \mapsto \theta_x(B) \text{ is measurable and } \theta_x(B) = \mathbb{E}(1_B|\mathcal{F})(x)\}$. From the definition of θ_x , $\mathcal{A} \subset \mathcal{C}$. By the properties of measures and by linearity of conditional expectation, we can see \mathcal{C} is an algebra. And by continuity of measures and conditional expectation, we can also check \mathcal{C} is a monotone class. So, by monotone class theorem, $\sigma(\mathcal{C}) = M(\mathcal{C}) = \mathcal{C}$. Thus, $\mathcal{B} = \sigma(\mathcal{A}) \subset \sigma(\mathcal{C}) = \mathcal{C}$. \square

Now, we extend $\{\theta_x\}_{x \in X_0}$ to the whole of X , by mapping $x \in X_0^c$ to a fixed probability measure. Hence, $\{\theta_x\}_{x \in X}$ is a kernel from (X, \mathcal{F}) to (X, \mathcal{B}) . And,

$$\int \theta_x(E) d\mu(x) = \int \mathbb{E}(1_E|\mathcal{F})(x) d\mu(x) = \mu(E).$$

Note,

$$\mathbb{E}(1_E|\mathcal{F})(x) = \int 1_E(t)d\theta_x(t).$$

So, by linearity and continuity of integral and conditional expectation, we can replace 1_E by any bounded \mathcal{B} -measurable functions $f : X \rightarrow \mathbb{R}$ to get (2.1).

Finally, the uniqueness of kernel follows from (2.1), by taking f to be characteristic functions. \square

Theorem 2.2.14 (Ergodic decomposition). *Let (X, \mathcal{B}, μ, T) be a standard measure preserving system, and $\mathcal{F} \subset \mathcal{B}$ the sub- σ -algebra of the T -invariant sets. If $\{\theta_x\}_{x \in X}$ is the disintegration of μ over \mathcal{F} , then θ_x is T -invariant and ergodic for μ -a.e.*

Proof.

Claim: $\theta'_x = \theta_x \circ T^{-1}$ is a disintegration of μ over \mathcal{F} .

Proof. Clearly, θ_x is a kernel. And,

$$\int \theta'_x(E)d\mu(x) = \int \theta_x(T^{-1}E)d\mu(x) = \mu(T^{-1}E) = \mu(E).$$

Lastly,

$$\int f(t)d\theta'_x(t) = \int f(t)d(\theta_x \circ T^{-1})(t) = \int f \circ T(t)d\theta_x(t) = \mathbb{E}(f \circ T|\mathcal{F})(x) = \mathbb{E}(f|\mathcal{F})(x) \text{ a.e.}$$

The second equality is due to a change of variable, and the last equality follows, since for any $A \in \mathcal{F}$,

$$\int_A \mathbb{E}(f \circ T|\mathcal{F})d\mu = \int_A f \circ T d\mu = \int_A f d\mu = \int_A \mathbb{E}(f|\mathcal{F})d\mu.$$

\square

Thus, by uniqueness of disintegration, $\theta_x \circ T^{-1} = \theta_x$, which means that θ_x is T -invariant. To show θ_x is ergodic, take $f : X \rightarrow \mathbb{R}$ a bounded, \mathcal{B} -measurable, T -invariant function. This means, \mathcal{F} -measurable, since for U open in \mathbb{C} ,

$$f \circ T = f \implies (f \circ T)^{-1}U = T^{-1}f^{-1}U = f^{-1}U \implies f^{-1}U \in \mathcal{F}.$$

So, $\mathbb{E}(f|\mathcal{F}) = f$ a.e. Now,

$$\begin{aligned} \text{Var}_x(f) &= \int \left[f(t) - \int f(s)d\theta_x(s) \right]^2 d\theta_x(t) = \int f(t)^2 d\theta_x(t) - \left(\int f(t)d\theta_x(t) \right)^2 \\ &= \mathbb{E}(f^2|\mathcal{F}) - \mathbb{E}(f|\mathcal{F})^2 = f^2 - f^2 = 0. \end{aligned}$$

Hence, the first integrand is 0. So, f is constant a.e. and θ_x is ergodic. \square

2.3 Ergodic Theorems

2.3.1 Mean ergodic theorem

Lemma 2.3.1. *If T is a contraction, then $T^*v = v \iff Tv = v$.*

Proof. If $T^*v = v$, then

$$\begin{aligned}\|Tv - v\|^2 &= \langle Tv - v, Tv - v \rangle \\ &= \|Tv\|^2 - \langle Tv, v \rangle - \langle v, Tv \rangle + \|v\|^2 \\ &\leq 2\|v\|^2 - 2\|v\|^2 = 0.\end{aligned}$$

So, $Tv = v$.

The other direction is similar. \square

Theorem 2.3.2 (Von Neumann's mean ergodic theorem for Hilbert space). *\mathcal{H} is a Hilbert space and T is a contraction [i.e. T is a bounded operator and $\|T\| \leq 1$]. Let $\mathcal{M} = \{v \in \mathcal{H} \mid Tv = v\}$ and $\pi : \mathcal{H} \rightarrow \mathcal{M}$ be the orthogonal projection. Then*

$$S_n(v) := \frac{1}{n} \sum_{k=0}^{n-1} T^k(v) \rightarrow \pi(v) \quad \forall v \in \mathcal{H}.$$

Proof. We write $\mathcal{H} = \mathcal{M} \oplus \mathcal{M}^\perp$.

Claim: Let $\mathcal{N} = \{v - Tv \mid v \in \mathcal{H}\}$, then $\mathcal{M}^\perp = \overline{\mathcal{N}}$.

Proof. Since $\overline{\mathcal{N}}, \mathcal{M}$ are closed, we can show $\mathcal{N}^\perp = \mathcal{M}$ (since $\overline{\mathcal{N}}^\perp = \mathcal{N}^\perp$). Let $v \in \mathcal{N}^\perp$, then

$$\langle v, w - Tw \rangle = 0 \implies \langle v, w \rangle = \langle v, Tw \rangle = \langle T^*v, w \rangle.$$

Since this is for every $w \in \mathcal{H}$, we get $T^*v = v$, and by the lemma, $Tv = v$. Hence, $v \in \mathcal{M}$.

Now, let $v \in \mathcal{M}$, then $Tv = v$ and $T^*v = v$. So,

$$\langle v, Tw - w \rangle = \langle T^*v, w \rangle - \langle v, w \rangle = 0$$

Hence, $v \in \mathcal{N}^\perp$. \square

By linearity of the operators, it is enough to show the statement holds on \mathcal{M} and \mathcal{M}^\perp separately. When $v \in \mathcal{M}$,

$$S_n(v) = v \rightarrow v$$

When $v \in \mathcal{M}^\perp = \overline{\mathcal{N}}$, by continuity, we only show when $v = w - Tw \in \mathcal{N}$

$$\begin{aligned}\|S_n(w - Tw)\| &= \left\| \frac{1}{n} \sum_{k=0}^{n-1} T^k(w - Tw) \right\| = \left\| \frac{1}{n} (w - T^n w) \right\| \\ &\leq \frac{1}{n} \|w\| + \|T^n w\| \leq \frac{2}{n} \|w\| \rightarrow 0 = \pi(v)\end{aligned}$$

\square

Let (X, \mathcal{B}, μ, T) be a MPS. We define an operator (for $1 \leq p \leq \infty$)

$$U_T : L^p(X, \mu) \rightarrow L^p(X, \mu), \quad U_T(f) = f \circ T.$$

Then U_T is an isometry, since (because T is μ -invariant)

$$\|U_T f\|_p = \left(\int |f(T(x))|^p d\mu(x) \right)^{1/p} = \left(\int |f(x)|^p d\mu(x) \right)^{1/p} = \|f\|_p, \quad (1 \leq p < \infty)$$

For $p = \infty$, it follows from the fact that $\|f\|_\infty = \lim_{p \rightarrow \infty} \|f\|_p$.

Corollary 2.3.3 (Mean ergodic theorem for dynamical systems). *Let (X, \mathcal{B}, μ, T) be a MPS and $M_T \subset \mathcal{B}$ be the σ -algebra of T -invariant subsets of \mathcal{B} , i.e. $M_T = \sigma\{E \in \mathcal{B} \mid T^{-1}E = E\}$, then for $f \in L^2(X, \mathcal{B}, \mu)$,*

$$S_n(U_T)(f) = \frac{1}{n} \sum_{k=0}^{n-1} f \circ T^k \rightarrow \mathbb{E}(f|M_T) \text{ in } \|\cdot\|_2.$$

If the system is ergodic, then

$$S_n(U_T)(f) \rightarrow \int f d\mu \text{ in } \|\cdot\|_2.$$

Proof. We use the above theorem on $\mathcal{H} = L^2(X, \mathcal{B}, \mu)$ and U_T as the contraction.

Let $\pi_T : L^2(X, \mathcal{B}, \mu) \rightarrow L^2(X, M_T, \mu)$ be the orthogonal projection, then $\pi_T(f) = \mathbb{E}(f|M_T)$. Indeed, $\pi_T(f)$ is clearly M_T -measurable and for $f \in L^2(X, \mathcal{B}, \mu)$, $E \in M_T$,

$$\int_E \pi_T(f) d\mu = \langle \pi_T(f), 1_E \rangle = \langle f, 1_E \rangle = \int_E f d\mu.$$

Finally, we show that $\mathcal{M} := \{f \in L^2(X, \mathcal{B}, \mu) \mid U_T(f) = f\} = L^2(X, M_T, \mu)$. First, take $1_A \in L^2(X, M_T, \mu)$, where $A \in M_T$. But, $U_T(1_A) = 1_{T^{-1}A} = 1_A$. So $1_A \in \mathcal{M}$. By linearity and continuity of U_T , $L^2(X, M_T, \mu) \subset \mathcal{M}$. For the reverse inclusion, we have to show that $f \in \mathcal{M}$ is M_T -measurable, or that for all $a \in \mathbb{R}$, $f^{-1}(-\infty, a)$ is T -invariant. This follows from $T^{-1}f^{-1}(-\infty, a) = (f \circ T)^{-1}(-\infty, a) = f^{-1}(-\infty, a)$.

When the system is ergodic, $M_T = \{\phi, X\}$. Since, $\mathbb{E}(f|M_T)$ is M_T -measurable, $\mathbb{E}(f|M_T)$ is a linear multiple of 1_X , $\mathbb{E}(f|M_T) = a1_X$. So, $a = \int \mathbb{E}(f|M_T) d\mu = \int f d\mu$. \square

We now look at the generalization of the theorem to Banach spaces. Let $(X, \|\cdot\|)$ be a Banach space, and $T : X \rightarrow X$ a continuous linear map such that $\sup_{n \in \mathbb{N}} \|T^n\| < \infty$. We call such operators, *power bounded operators*. Note, for such T ,

$$\lim_{n \rightarrow \infty} \frac{T^n(x)}{n} = 0, \text{ and } \sup_{n \in \mathbb{N}} \|S_n\| < \infty.$$

We define the average

$$S_n(x) = \frac{1}{n} \sum_{k=0}^{n-1} T^k(x).$$

Consider the following subspaces,

$$\begin{aligned} N &= \{x \in X \mid T(x) = x\} = \ker(I - T), \\ M &= \{x \in X \mid \lim_{n \rightarrow \infty} S_n(x) \text{ exists}\}. \end{aligned}$$

And define the map,

$$P : M \rightarrow X, \quad P(x) = \lim_{n \rightarrow \infty} S_n(x).$$

Lemma 2.3.4. *P has the following properties,*

- (i) $P \in \mathcal{B}(X)$
- (ii) $P(x) = PT(x) = TP(x)$.
- (iii) $x \in \overline{\text{Im}(I - T)} \implies P(x) = 0$

Proof. (i)

$$\|P(x)\| = \lim_{n \rightarrow \infty} \|S_n\| \|x\| \implies \|P\| < \infty,$$

since $\sup_n \|S_n\| < \infty$.

(ii)

$$S_n(Tx) = \frac{1}{n} \sum_{k=0}^{n-1} T^k(Tx) = \frac{1}{n} \sum_{k=0}^{n-1} T^k(x) + \frac{T^n x - x}{n} \implies PT(x) = P(x),$$

since $\frac{T^n(x)}{n} \rightarrow 0$. And

$$S_n(Tx) = \frac{1}{n} \sum_{k=0}^{n-1} T^k(Tx) = T \left(\frac{1}{n} \sum_{k=0}^{n-1} T^k(x) \right) \implies PT(x) = TP(x),$$

by linearity and continuity of T .

(iii) If $x = y - Ty \in \text{Im}(I - T)$, for some $y \in X$. Then, $P(x) = P(y) - PT(y) = 0$. By continuity of P , we can extend to x in the closure. □

Theorem 2.3.5. *The following hold,*

- (i) M is a T -invariant closed subspace of X .
- (ii) $P : M \rightarrow N$ is a projection, i.e., $P^2 = P$, and $M = N \oplus \overline{\text{Im}(I - T)}$.

Proof. M is T -invariant: We have to show if $x \in M$, then $Tx \in M$. This follows from $PT(x) = P(x)$.

M is closed: Consider a sequence $x_n \in M$ such that $x_n \rightarrow x$. To show $\lim S_n(x)$ exists, we show $S_n(x)$ is Cauchy. Given $\varepsilon > 0$, choose k such that $\|x - x_k\| < \varepsilon$. Since $S_n(x_k)$ is Cauchy, choose m, n such that $\|(S_n - S_m)(x_k)\| < \varepsilon$. Then,

$$\begin{aligned} \|(S_n - S_m)(x)\| &\leq \|S_n(x - x_k)\| + \|(S_n - S_m)(x_k)\| + \|S_m(x_k - x)\| \\ &\leq (2 \sup_{n \in \mathbb{N}} \|S_n\| + 1)\varepsilon \rightarrow 0. \end{aligned}$$

$P^2 = P$: $\text{Im}(T) \subset N$ follows from the fact that $TP = P$. Since elements of N are T fixed points and hence S_n fixed points, this also shows that $P^2 = P$.

$M = N \oplus \overline{\text{Im}(I - T)}$: For $x \in M$, $x = P(x) + (I - P)(x)$. Since, $(I - S_n)(x) \in \text{Im}(I - T)$, we have $(I - P)(x) \in \overline{\text{Im}(I - T)}$. And if $x \in N \cap \overline{\text{Im}(I - T)}$, then from the lemma, $P(x) = x = 0$. □

Theorem 2.3.6. *For $x, y \in X$, the following are equivalent,*

- (i) $\lim_{n \rightarrow \infty} S_n x = y$.
- (ii) there is a subsequence (n_k) such that $S_{n_k}(x) \rightarrow y$.
- (iii) $y \in N \cap \overline{\text{conv}^{weak}\{T^n x \mid n \geq 0\}}$.
- (iv) $y \in N \cap \overline{\text{conv}}^{\|\cdot\|}\{T^n x \mid n \geq 0\}$.

Proof. (i) \implies (ii): This is due to the fact that norm-convergent sequences are also weakly-convergent.

(iii) \implies (iv): This is a consequence of convex sets having same closure in norm topology and weak topology.

(ii) \implies (iii): Clearly, $y \in \overline{\text{conv}}^{\text{weak}}\{T^n x \mid n \geq 0\}$. Showing $y \in N$ is similar to showing $TP = P$, but using weak-continuity of T instead.

(iv) \implies (i): Since $y \in \overline{\text{conv}}^{\|\cdot\|}\{T^n x \mid n \geq 0\}$, there is a sequence $\{y_n\} \in \text{conv}^{\|\cdot\|}\{T^n x \mid n \geq 0\}$, and $y_n = \sum_{i=1}^{k_n} c_{n_i} T^{n_i}(x)$. Now,

$$\begin{aligned} y_n &= x - (x - y_n) = x - \left(\sum_{i=1}^{k_n} c_{n_i} \right) (I - T^{n_i})(x) \\ &= x - (I - T)(x'_n), \end{aligned}$$

for some $x'_n \in X$, since $I - T^{n_i} = (I - T)(I + T + T^2 + \cdots + T^{n_i-1})$. Hence, $y_n = x - x'_n$, where $x'_n \in \text{Im}(I - T)$. Taking limits, $y = x - x'$, where $x' \in \text{Im}(I - T)$. Thus, $\lim S_n y = y = \lim S_n x$, as $\lim S_n x' = 0$, by the lemma. \square

Theorem 2.3.7 (Mean ergodic theorem for Banach spaces). *Let X be a reflexive Banach space, and $T \in \mathcal{B}(X)$ such that $\sup_{n \in \mathbb{N}} \|T^n\| < \infty$, then*

$$\lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} T^k(x) \text{ exists for all } x \in X.$$

Proof. We know, in a reflexive Banach space, every bounded sequence has a weakly convergent subsequence. The result follows, by the previous theorem. \square

2.3.2 Pointwise ergodic theorem

(X, \mathcal{B}, μ, T) is a measure-preserving system. Define the average operator on $L^1(\mu)$,

$$S_n f = \frac{1}{n} \sum_{k=0}^{n-1} f \circ T^k.$$

Let \mathcal{F} be the σ -algebra of T -invariant sets.

Theorem 2.3.8 (Birkhoff's pointwise ergodic theorem). *For $f \in L^1(\mu)$, the following holds,*

$$S_n f(x) \rightarrow \mathbb{E}(f|\mathcal{F})(x) \text{ a.e.} \quad (2.2)$$

If the system is ergodic, then $S_n f(x) \rightarrow \int f d\mu$ a.e.

For the proof, we first find a dense set of L^1 where the statement holds and then attempt to extend to L^1 .

Lemma 2.3.9. *There is a dense set $S \subset L^1$ such that for $f \in L^1$, (2.2) holds.*

Proof. Since L^2 is dense in L^1 , it is sufficient to find a dense set of L^2 . Define $S_1 = \{f \in L^2 \mid fT = f\}$ and $S_2 = \{g - gT \mid g \in L^\infty\}$. Clearly, (2.2) holds on S_1 . On S_2 , $\frac{1}{n} S_n(g - gT) = \frac{1}{n}(g - gT^n) \rightarrow 0$, and $\int_F \mathbb{E}(g - gT|\mathcal{F}) = \int_F g - gT = 0$.

Claim: $\overline{S_2} = \overline{\{g - gT \mid g \in L^2\}}$ in L^2 .

Proof. Since $L^\infty \subset L^2$, we have $\overline{S_2} \subset \overline{\{g - gT \mid g \in L^2\}}$. Now, let $g \in L^2$. By density of L^∞ in L^2 , there is a $g' \in L^\infty$ such that $\|g - g'\|_2 < \varepsilon/2$. Then, $\|g - gT - (g' - g'T)\| < \|g - gT\| + \|g' - g'T\| < \varepsilon$. Hence, $g - gT \in \overline{S_2}$ and $\overline{\{g - gT \mid g \in L^2\}} \subset \overline{S_2}$. \square

We saw in the proof of the mean ergodic theorem that, $L^2 = S_1 \oplus \overline{S_2}$. But, $S_1 \oplus \overline{S_2} \subset \overline{S_1 \oplus S_2}$. Thus, $S_1 \oplus S_2$ is the required dense set. \square

To extend to all of L^1 , we need the maximal inequality. But before that, we look at its discrete version.

If $\hat{f} : \mathbb{N} \rightarrow [0, \infty)$, define the average over $I \subset \mathbb{N}$,

$$S_I f = \frac{1}{|I|} \sum_{i \in I} \hat{f}(i).$$

Theorem 2.3.10 (Discrete maximal inequality). *Let $\hat{f} : \mathbb{N} \rightarrow [0, \infty)$, $I \subset \mathbb{N}$ be finite interval and $J \subset I$ be subset such that $\forall j \in J$, $I_j \subset I$ be sub-intervals of I with left-endpoint j . If $S_{I_j} \hat{f} > t$, then*

$$S_I \hat{f} > t \frac{|J|}{|I|}.$$

Proof. First, if I_j are disjoint, then $\{I_j\}_j, \{I \setminus \bigcup_j I_j\}$ partition I . So,

$$S_I \hat{f} = \frac{|I \setminus \bigcup_j I_j|}{|I|} S_{I \setminus \bigcup_j I_j} \hat{f} + \sum \frac{|I_j|}{|I|} S_{I_j} \hat{f} \geq \frac{|\bigcup_j I_j|}{|I|} S_{I_j} \hat{f} \geq \frac{|J|}{|I|} t.$$

If I_j are not disjoint, we find a sub-collection $J_0 \subset J$ such that $\{I_j\}_{j \in J_0}$ are disjoint and $J \subset \bigcup_{j \in J_0} I_j$. Then the above inequalities will still go through, with union over J_0 instead of J . Let $j_1 = \min J$ and define $j_k = \min(J \setminus \bigcup_{i=1}^{k-1} I_{j_i})$. \square

Theorem 2.3.11 (Maximal inequality). *$f \in L^1$, $f \geq 0$, then $\forall t > 0$,*

$$\mu\{x \in X \mid \sup_n S_n f(x) > t\} \leq \frac{1}{t} \int f d\mu.$$

Proof. Let $A = \{\sup_n S_n f > t\}$. Fix $x \in X$, define $\hat{f}_x : \mathbb{N} \rightarrow [0, \infty)$, $\hat{f}_x(i) = f(T^i x)$. Now, we try to use the discrete maximal inequality on \hat{f}_x . If $T^j x \in A$, then $S_{n_j} T^j x > t$ for some n_j . Let $I_j = [j, j + n_j - 1]$, then

$$S_{I_j} \hat{f}_x = \frac{1}{n_j} \sum_{i=j}^{j+n_j-1} f(T^i x) = \frac{1}{n_j} \sum_{i=0}^{n_j-1} f(T^i(T^j x)) = S_{n_j} T^j x > t.$$

Consider $I = [0, M - 1]$, $J = \{j \in I \mid T^j x \in A, I_j \subset I\}$. By the discrete maximal inequality on $\{I_j\}_{j \in J}$,

$$S_I \hat{f}_x > t \frac{|J|}{M}.$$

Let $A_R = \{\sup_{0 \leq n \leq R} S_n f > t\}$, then $J \supset \{j \in [0, M - R] =: I' \mid T^j x \in A_R\}$ and

$$|J| \geq \sum_{j=0}^{M-R} 1_{A_R}(T^j x) = \sum_{j=0}^{M-R} \hat{1}_{A_R}(j) = (M - R) S_{I_R} \hat{1}_{A_R}.$$

Thus, $S_I \hat{f}_x > t \frac{|J|}{M} \geq t \frac{M-R}{M} S_{I_R} \hat{1}_{A_R}$. Now,

$$\begin{aligned} \int f d\mu &= \frac{1}{M} \sum_{i=0}^{M-1} \int f T^i(x) d\mu(x) = \int S_I \hat{f}_x d\mu(x) > t \frac{M-R}{M} \int S_{I_R} \hat{1}_{A_R, x} d\mu(x) \\ &= t \left(1 - \frac{R}{M}\right) \mu(A_R) \xrightarrow{M \rightarrow \infty} t \mu(A_R) \xrightarrow{R \rightarrow \infty} t \mu(A) \end{aligned}$$

The last equality follows from,

$$\int S_{I'} \hat{1}_{A_R, x} d\mu(x) = \int \frac{1}{M-R} \sum_{i=0}^{M-R} \hat{1}_{A_R, x}(i) d\mu(x) = \frac{1}{M-R} \sum_{i=0}^{M-R} \int 1_{A_R}(T^i x) d\mu(x) = \int 1_{A_R} d\mu.$$

□

Now we can complete the proof of the pointwise ergodic theorem.

Proof of pointwise ergodic theorem. We write $S(f) = \mathbb{E}(f|\mathcal{F})$, for convenience. Let $f \in L^1$, and take $g \in S$. Then, a.e. $x \in X$,

$$|S_n f - S f| \leq |S_n f - S_n g| + |S_n g - S f| \leq S_n |f - g| + |S_n g - S f|.$$

But, $S_n g \rightarrow S g$ a.e., implies $|S_n g - S f| \rightarrow |S g - S f| \leq S|g - f|$, by property of conditional expectation. So,

$$\limsup_{n \rightarrow \infty} |S_n f - S f| \leq \limsup_{n \rightarrow \infty} S_n |f - g| + S|g - f|.$$

Note, if $\limsup_{n \rightarrow \infty} |S_n f - S f| > \varepsilon$, then either $\limsup_{n \rightarrow \infty} S_n |f - g| > \varepsilon/2$ or $S|g - f| > \varepsilon/2$. Therefore,

$$\mu\{\limsup_{n \rightarrow \infty} |S_n f - S f| > \varepsilon\} \leq \mu\{\limsup_{n \rightarrow \infty} S_n |f - g| > \frac{\varepsilon}{2}\} + \mu\{S|g - f| > \frac{\varepsilon}{2}\} \leq \frac{2}{\varepsilon} \|f - g\|_1 + \frac{2}{\varepsilon} \|f - g\|_1.$$

The first term comes from the maximal inequality and the second term from Markov's inequality. Since the right hand side can be made arbitrarily small, we have $\mu\{\limsup_{n \rightarrow \infty} |S_n f - S f| > \varepsilon\} = 0$. Thus, $\limsup_{n \rightarrow \infty} |S_n f - S f| = 0$ a.e., and $S_n f \rightarrow S f$ a.e. □

2.4 Mixing

Definition 2.4.1 (Strongly mixing). (X, \mathcal{B}, μ, T) is *strong mixing* if for all $A, B \in \mathcal{B}$,

$$\lim_{n \rightarrow \infty} \mu(A \cap T^{-n} B) = \mu(A)\mu(B).$$

It is clear that strongly ergodic systems are ergodic.

Theorem 2.4.2. (X, \mathcal{B}, μ, T) is *strongly mixing* if and only if for every $f, g \in L^2$,

$$\lim_{n \rightarrow \infty} \int f \cdot (g T^n) d\mu = \int f d\mu \cdot \int g d\mu.$$

Furthermore, the above is true for every $f, g \in L^2$ if it holds for a dense set of L^2 .

Proof. We first see how the second statement implies the first. For $A, B \in \mathcal{B}$, take $f = 1_A, g = 1_B$. Then $1_A \cdot 1_B(T^n) = 1_A \cdot 1_{T^{-n} B} = 1_{A \cap T^{-n} B}$. So, the statement holds for simple functions, and by the density of simple functions in L^2 , statement holds for the entire L^2 .

Now, suppose $S \subset L^2$ is a dense subset such the limit holds on S . For $f, g \in L^2$, let $f', g' \in S$, such that $\|f' - f\| < \varepsilon, \|g' - g\| < \varepsilon$. Then,

$$\begin{aligned} \left| \int f \cdot g T^n - \int f \cdot \int g \right| &\leq \left| \int (f - f' + f') \cdot (g - g' + g') T^n - \int f \cdot \int g \right| \\ &\leq \left| \int (f - f') \cdot (g - g') T^n \right| + \left| \int (f - f') \cdot g' T^n \right| + \left| \int f' \cdot (g - g') T^n \right| \\ &\quad + \left| \int f' \cdot g' T^n - \int f' \cdot \int g' \right| + \left| \int f' \cdot \int g' - \int f \cdot \int g' \right| \\ &\quad + \left| \int f \cdot \int g' - \int f \cdot \int g \right| \\ &\leq \varepsilon^2 + \varepsilon \|g'\| + \varepsilon \|f'\| + \varepsilon + \varepsilon \|g'\| + \varepsilon \|f\| \rightarrow 0. \end{aligned}$$

The last inequality is follows applications of Cauchy-Schwarz inequality. \square

The above technique will be used repeatedly to extend similar statements from a dense subset to the entire space.

We also have a similar characterization for ergodicity.

Theorem 2.4.3. (X, \mathcal{B}, μ, T) is ergodic if and only if $A, B \in \mathcal{B}$,

$$\lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} \mu(A \cap T^{-k}B) = \mu(A)\mu(B).$$

Proof. \Leftarrow : If $\mu(A), \mu(B) > 0$, then $\mu(A \cap T^{-n}B) > 0$, for infinitely many n . If not, then the limit would go to 0, which is not possible.

\Rightarrow : For $A, B \in \mathcal{B}$,

$$\begin{aligned} \lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} \mu(A \cap T^{-k}B) &= \lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} \int 1_A \cdot 1_B(T^k) = \lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} \langle 1_A, 1_B(T^k) \rangle \\ &= \left\langle 1_A, \lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} 1_B(T^k) \right\rangle = \left\langle 1_A, \lim_{n \rightarrow \infty} S_n 1_B \right\rangle \\ &= \left\langle 1_A, \int 1_B \right\rangle = \mu(A)\mu(B). \end{aligned}$$

The penultimate equality is due to mean ergodic theorem. \square

As before, we can similarly extend to integrals of functions instead of measures of sets.

Definition 2.4.4 (Weakly mixing). (X, \mathcal{B}, μ, T) is *weakly mixing* if for all $A, B \in \mathcal{B}$,

$$\lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} |\mu(A \cap T^{-k}B) - \mu(A)\mu(B)| = 0.$$

We can again extend this to functions instead of sets: for all $f, g \in L^2$,

$$\lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} \left| \int f \cdot g T^k - \int f \int g \right| = 0.$$

We require the concept of density, before studying weakly mixing further.

Definition 2.4.5 (Upper density). The density of subset $I \subset \mathbb{N}$ is

$$d(I) = \lim_{n \rightarrow \infty} \frac{|I \cap \{1, \dots, n\}|}{n}.$$

Definition 2.4.6 (Convergence in density). a_n is said to *converge in density* to a , $d\text{-}\lim_{n \rightarrow \infty} a_n = a$ or $a_n \xrightarrow{d} a$, if

$$d(\{n \mid |a_n - a| > \varepsilon\}) = 0 \text{ for all } \varepsilon > 0.$$

The definition of weakly mixing can now be reformulated as $\mu(A \cap T^{-n}B) \xrightarrow{d} \mu(A)\mu(B)$.

Notice, usual convergence implies convergence in density, because the above set will always have finite elements and hence, zero density. So, we have strong mixing implies weak mixing.

But the converse is not true. The sequence $a_n = 1$, when n is prime and 0 otherwise, will converge to 0, as the primes have 0 density, but does not converge in norm.

Lemma 2.4.7. *For bounded sequences a_n , the following are equivalent,*

- (i) $d\text{-}\lim_{n \rightarrow \infty} a_n = a$.
- (ii) *there is a subset $J \subset \mathbb{N}$ with $d(J) = 0$ such that $\lim_{n \rightarrow \infty, n \notin J} a_n = a$.*
- (iii) $\lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} |a_k - a| = 0$.
- (iv) $\lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} (a_k - a)^2 = 0$.

Proof. (ii) \implies (i): Follows from the fact that, $\{i \in [1, n] \mid |a_i - a| \geq \varepsilon\} \subset \{i \in [1, n] \cap J^c \mid |a_i - a| \geq \varepsilon\} \cup \{i \in [1, n] \cap J \mid |a_i - a| \geq \varepsilon\}$. The first set is finite and the second has density zero.

(i) \implies (iii): For $\varepsilon > 0$, the set $\{n \mid |a_n - a| \geq \varepsilon\}$ has density zero. And since the sequences are bounded, $|a_n - a| < K$ for some $K > 0$. Thus,

$$\frac{1}{n} \sum_{k=0}^{n-1} |a_k - a| < \varepsilon + \frac{K}{n} \{i \in [1, n] \mid |a_i - a| \geq \varepsilon\} \rightarrow \varepsilon.$$

(iii) \implies (ii): Define $J_k = \{n \mid |a_n - a| \geq 1/k\}$. Then $J_k \subset J_{k+1}$ and each J_k has density zero, since $\frac{1}{n} \sum |a_n - a| \geq \frac{1}{kn} |J_k \cap [1, n]|$. This also means that, there is a sequence $\{l_k\}$ such that $|J_k \cap [1, n]|/n < 1/k \forall n \geq l_k$. Now, define $J = \bigcup_k (J_k \cap [l_k, l_{k+1}))$, then J also has density zero. Indeed, $J \cap [1, n] \subset J_k \cap [1, n]$, where $l_k \leq n < l_{k+1}$. Finally, if $n > l_k$ and $n \notin J$, then $n \notin J_k$, and hence $|a_n - a| < 1/k$. □

Now, we can see weakly mixing implies ergodic. If it was not ergodic, then the series (ii) cannot converge to zero, by the previous theorem on ergodicity. Using (ii), we can also easily show convergence in density respects sums and products.

Corollary 2.4.8. *Strong mixing implies weak mixing, and weak mixing implies ergodicity.*

Theorem 2.4.9. *(X, \mathcal{B}, μ, T) is weakly mixing if and only if $(X \times X, \mathcal{B} \times \mathcal{B}, \mu \times \mu, T \times T)$ is ergodic.*

Proof. \Leftarrow : We can see the condition implies ergodicity. If T is not ergodic $T \times T$ cannot be ergodic, since if f is a non-constant T -invariant measurable function, then (f, f) is a non-constant $T \times T$ -invariant measurable function.

Let $f, g \in L^2(\mu)$, and define $\tilde{f}, \tilde{g} \in L^2(\mu \times \mu)$ by $\tilde{f}(x, y) = f(x)f(y)$, $\tilde{g}(x, y) = g(x)g(y)$. Since $T \times T$ is ergodic,

$$\frac{1}{n} \sum_{k=0}^{n-1} \int \tilde{f} \cdot \tilde{g}(T \times T)^k \rightarrow \int \tilde{f} \int \tilde{g}.$$

By Fubini's theorem, the left integral is $(\int f \cdot g T^k)^2$ and the right integral is $(\int f \int g)^2$. Thus, $\frac{1}{n} \sum_{k=0}^{n-1} (\int f \cdot g T^k)^2 \rightarrow (\int f \int g)^2$. Note, since T is ergodic, $\frac{1}{n} \sum_{k=0}^{n-1} \int f \cdot g T^k \rightarrow \int f \int g$. Finally,

$$\begin{aligned} \lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} \left(\int f \cdot g T^k - \int f \cdot \int g \right)^2 &= \lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} \left[\left(\int f \cdot g T^k \right)^2 + \left(\int f \int g \right)^2 \right. \\ &\quad \left. - 2 \left(\int f \cdot g T^k \right) \left(\int f \int g \right) \right] = 0. \end{aligned}$$

\implies : We prove for $A_1 \times A_2, B_1 \times B_2 \in \mathcal{B} \times \mathcal{B}$.

$$\begin{aligned} (\mu \times \mu) ((A_1 \times A_2) \cap (T \times T)^{-k}(B_1 \times B_2)) &= (\mu \times \mu) ((A_1 \cap T^{-k}B_1) \cap (A_2 \cap T^{-k}B_2)) \\ &= \mu(A_1 \cap T^{-k}B_1) \mu(A_2 \cap T^{-k}B_2) \xrightarrow{d} \mu(A_1) \mu(A_2) \mu(B_1) \mu(B_2) = \mu(A_1 \times A_2) \mu(B_1 \times B_2). \end{aligned}$$

We know, functions of the form $1_{A_1}(x) \cdot 1_{A_2}(y) = 1_{A_1 \times A_2}(x, y)$, $A_1, A_2 \in \mathcal{B}$ form an orthonormal basis for $L^2(\mu \times \mu)$. Thus, by the density argument, the above convergence can be extended to all sets in the product σ -algebra, and we have $T \times T$ is ergodic. \square

Now, we discuss some multiplier properties of weak mixing systems.

Theorem 2.4.10. *(X, \mathcal{B}, μ, T) is weak mixing if and only if $X \times Y$ is ergodic for every ergodic system (Y, \mathcal{C}, ν, S) .*

Proof. \Leftarrow : Take $Y = \{*\}$, then $X \times \{*\} \cong X$ is ergodic. So, $X \times X$ is ergodic, and X is weak mixing

\implies : Again by density argument, it is enough to consider cylindrical sets $A \times C, B \times D \in \mathcal{B} \times \mathcal{C}$,

$$\begin{aligned} \frac{1}{n} \sum_{k=0}^{n-1} (\mu \times \nu) [(T \times S)^{-k}(A \times C) \cap (B \times D)] &= \frac{1}{n} \sum_{k=0}^{n-1} \mu(T^{-k}A \cap B) \nu(S^{-k}C \cap D) \\ &= \frac{1}{n} \sum_{k=0}^{n-1} [\{\mu(T^{-k}A \cap B) - \mu(A)\mu(B)\} \nu(S^{-k}C \cap D) + \mu(A)\mu(B) \nu(S^{-k}C \cap D)] \\ &\rightarrow \mu(A)\mu(B)\nu(C)\nu(D), \text{ as } n \rightarrow \infty. \end{aligned}$$

The first term goes to 0 and second term converges, due to ergodicity of μ and ν respectively. \square

Corollary 2.4.11. (i) X_1, X_2 is weak mixing $\implies X_1 \times X_2$ is weak mixing.

(ii) T is weak mixing $\implies T^n$ is weak mixing for all $n \in \mathbb{N}$

(iii) if T is invertible, T is weak mixing $\iff T^{-1}$ is weak mixing

Proof. (i) Take ergodic system Y , then $X_2 \times Y$ is ergodic, which implies $X_1 \times X_2 \times Y$ is also ergodic and hence, $X_1 \times X_2$ is weak mixing.

(ii) We first show, T is weak mixing $\implies T^n$ is ergodic for all $n \in \mathbb{N}$.

Assume T^n is not ergodic, and take a non-constant a.e., T^n -invariant function $f \in L^2(\mu)$. Define a measure-preserving system, (Y, \mathcal{C}, ν, S) where $Y = \{0, 1, \dots, n-1\}$, $\mathcal{C} = \mathcal{P}(Y)$, $\nu(A) = \frac{|A|}{n}$ and $S : Y \rightarrow Y$, $S(i) = i + 1 \pmod n$. It is easy to see that S is an ergodic system. So, $T^n \times S$ is ergodic. Define $F \in L^2(\mu \times \nu)$, $F(x, i) = f(T^{n-i}x)$. Then, F is $T \times S$ -invariant: $F((T \times S)(x, i)) = F(Tx, i + 1) = f(T^{n-i}x) = F(x, i)$. As, F is non-constant a.e., this contradicts the ergodicity of $T \times S$. Thus, T^n is ergodic.

Now, T is weak mixing $\implies T \times T$ is weak mixing $\implies T^n \times T^n$ is ergodic $\implies T^n$ is weak mixing.

(iii) Follows from the fact that T is ergodic if and only if T^{-1} is ergodic, which implies $T \times T$ is ergodic if and only if $T^{-1} \times T^{-1}$ is ergodic. \square

We now look at the relationship between properties of a measure-preserving system (X, \mathcal{B}, μ, T) and the point-spectrum of U_T . Recall, $\lambda \in \mathbb{C}$ is an *eigenvalue* of U_T if $U_T f = \lambda f$, for some $0 \neq f \in L^2(\mu)$, and f is called *eigenfunction*. The vector space of $f \in L^2(\mu)$ satisfying the equation is called the *eigenspace*, E_λ and the *point-spectrum* of U_T is defined by $\sigma_p(U_T) = \{\lambda \in \mathbb{C} \mid U_T f = \lambda f, f \in L^2(\mu)\}$. An eigenvalue is said to be *simple* if the corresponding eigenspace is one-dimensional.

Theorem 2.4.12. *If T is ergodic, then σ_p is a subgroup of S^1 , and every eigenvalue of U_T is simple. T is ergodic if and only if 1 is a simple eigenvalue of U_T .*

Proof. $\sigma_p \subset S^1$: If $\lambda \in \sigma_p$, then $U_T f = \lambda f$ for some non-zero $f \in L^2$. So, $\|f \circ T\| = \|f\| = |\lambda| \|f\|$, and $\lambda \in S^1$.

σ_p is closed under product: $\lambda_1, \lambda_2 \in \sigma_p$, and $U_T f = \lambda_1 f$ and $U_T g = \lambda_2 g$ for non-zero $f, g \in L^2$. Then $U_T(f \cdot g) = \lambda_1 \lambda_2 (f \cdot g)$, and $\lambda_1 \lambda_2 \in \sigma_p$.

σ_p is closed under inverse: If $U_T f = \lambda f$, then $|f \circ T| = |f|$. So, $|f|$ is T -invariant and hence, non-zero constant a.e. Thus, since μ is probability measure, $\|f^{-1}\|_2 < \infty$ and $f^{-1} \in L^2$. Now, $U_T(f \cdot f^{-1}) = U_T(1) = U_T(f)U_T(f^{-1}) \implies U_T(f^{-1}) = \lambda^{-1} f^{-1}$. Hence, $\lambda^{-1} \in \sigma_p$.

eigenvalues are simple: Consider $f, g \in E_\lambda$, then $U_T(f/g) = (\lambda f)/(\lambda g) = f/g$ and $f/g = c$ for some $c \in \mathbb{C}$. Thus, E_λ is one-dimensional.

\iff : Follows from the observation that 1 is a simple eigenvalue of U_T is equivalent to the statement that T -invariant functions are constant a.e. □

Theorem 2.4.13. *If T is ergodic, then it is weak mixing if and only if $\sigma_p = \{1\}$.*

Proof. \implies : Let $1 \neq \lambda \in \sigma_p$ and non-zero $f \in L^2$ such that $U_T f = \lambda f$. Then $\int f = \int U_T f = \lambda \int f$. As $\lambda \neq 1$, $\int f = 0$. Now,

$$\frac{1}{n} \sum_{k=0}^{n-1} \left| \int f T^k \cdot g - \int f \int g \right| = \frac{1}{n} \sum_{k=0}^{n-1} |\langle U_T^k f, \bar{g} \rangle| = \langle f, \bar{g} \rangle.$$

The last equality is due to $|\lambda| = 1$. Since, T is weak mixing, the above quantity goes to 0 as $n \rightarrow \infty$. Thus, we get $\langle f, \bar{g} \rangle = 0$ for all $g \in L^2$. So, $f = 0$, which is a contradiction.

\impliedby : We show $T \times T$ is ergodic. Let $k \in L^2(\mu \times \mu)$ be a $T \times T$ -invariant function. Define $k^*(y, x) = \overline{k(x, y)}$. Then k^* is also $T \times T$ -invariant. We can assume $k = k^*$, since we can write any k as sum of such functions, $k = \frac{1}{2}(k + k^*) + \frac{1}{2}i(k - k^*)$. Define $A : L^2(\mu) \rightarrow L^2(\mu)$, $Af(x) = \int k(x, y)f(y)d\mu(y)$.

A is a well-defined bounded operator:

$$\|Af\|_2^2 \leq \int \left(\int |k(x, y)f(y)|d\mu(y) \right)^2 d\mu(x) \leq \|k\|_2^2 \|f\|_2^2.$$

The last inequality is due to Cauchy-Schwarz inequality. Thus, $Af \in L^2$ and A is bounded. A is self-adjoint:

$$\langle f, A^*g \rangle = \langle Af, g \rangle = \iint k(x, y)f(y)\bar{g}(x)dydx = \iint \overline{F^*(y, x)g(x)}dx f(y)dy = \langle f, Ag \rangle$$

A is compact: We show that A is a Hilbert-Schmidt operator and use the fact that Hilbert-Schmidt operators are compact.

$$AU_T = U_TA:$$

$$\begin{aligned} AU_T(f)(x) &= Af(Tx) = \int F(x, y)f(Ty)dy = \int F(Tx, Ty)f(Ty)dy \\ &= \int F(Tx, y)f(y)dy = U_TA(f)(x) \end{aligned}$$

Consider the eigenspaces E_λ corresponding to non-zero eigenvalues λ of A . Then $A|_{E_\lambda} = \lambda \cdot \text{id}$ is compact. We know, $\text{id} : E_\lambda \rightarrow E_\lambda$ is compact if and only if E_λ is finite-dimensional. We also have $U_T(E_\lambda) \subset E_\lambda$: if $f \in E_\lambda$, then $AU_T(f) = U_TA(f) = \lambda U_T(f)$, and $U_T(f) \in E_\lambda$. Thus we get $U_T : E_\lambda \rightarrow E_\lambda$ is an isometry between finite-dimensional space, which makes it unitary (as isometries are injective) and hence, diagonalizable. But, $\sigma_p(U_T) = \{1\}$ and by previous theorem, 1 is a simple eigenvalue. Thus, E_λ are one-dimensional and with A and U_T commuting, every eigenfunction of A is also an eigenfunction of U_T . But, 1 is the only eigenfunction of U_T and hence of A . Spectral theorem for self-adjoint compact operator states that $Af = \sum_n \lambda_n \langle f, e_n \rangle e_n$, where λ_n are eigenvalues with eigenfunctions e_n . So, $Af = \lambda \langle f, 1 \rangle 1$, where λ is the eigenvalue of 1. Finally,

$$\begin{aligned} Af(x) &= \lambda \int f(y)dy = \iint k(x, y)f(y)dy \\ \implies \iint [k(x, y) - \lambda]f(y)dy &= 0 \\ \implies \iint [k(x, y) - \lambda]f(y)g(x) &= 0, \text{ for all } f, g \in L^2 \end{aligned}$$

Hence, k is constant a.e., and T is weak-mixing. □

Chapter 3

Multiple Recurrence Theorem

3.1 Szemerédi's theorem

An interesting generalization of the Poincaré recurrence is the multiple recurrence theorem, which states any positive measure set has a positive measure subset that returns to it in an arithmetic progression of arbitrary length. It has many applications, an important one being that it connects ergodic theory and combinatorial number theory.

Theorem 3.1.1 (Multiple Recurrence Theorem). *For a measure-preserving system (X, \mathcal{B}, μ, T) , and a set $A \in \mathcal{B}$, with $\mu(A) > 0$, and for any $k \in \mathbb{N}$, there is some $n \geq 1$ such that*

$$\mu(A \cap T^{-n}A \cap T^{-2n}A \cap \dots T^{-kn}A) > 0$$

When $k = 1$, this reduces to Poincaré recurrence.

We now use this theorem to prove Szemerédi's theorem, which states that any set of integers with positive density contains arithmetic progressions of arbitrary length.

Definition 3.1.2 (Upper density). The *upper density* of a set $A \subset \mathbb{Z}$ is defined as

$$d(A) = \limsup_{N \rightarrow \infty} \frac{1}{2N+1} |S \cap \{-n, -n+1, \dots, n-1, n\}|$$

Theorem 3.1.3 (Szemerédi's Theorem). *If $A \subset \mathbb{Z}$ has positive upper density, then A contains arithmetic progressions of arbitrary length, that is, for all $k \in \mathbb{N}$, there exists $a, b \in \mathbb{Z}$, $b \neq 0$ such that $a, a+b, a+2b, \dots, a+kb \in A$.*

Proof. Define a MPS in the following way. $X = \{0, 1\}^{\mathbb{Z}}$ with the Bernoulli shift

$$\sigma : X \rightarrow X, (\sigma(x))_n = x_{n-1}.$$

For a set $S \subset \mathbb{Z}$, define $x \in X$ by,

$$x_n = \begin{cases} 1 & \text{if } n \in S, \\ 0 & \text{if } n \notin S. \end{cases}$$

For $n \in \mathbb{N}$, define the probability measure on X ,

$$\mu_n = \frac{1}{2n+1} \sum_{k=-n}^n \delta_{\sigma^k(x)},$$

where δ_x is the Dirac delta measure on $x \in X$. Now, note the set of probability measures on X , $\mathcal{P}(X)$ is weak*-compact (by using Riesz theorem and Banach-Alaoglu theorem). Recall, if X is

metrizable and compact, then $C(X)$ is separable. Then, $\mathcal{P}(X)$ is metrizable with the metric: let $\{f_i\}$ be a countable dense subset of $C(X)$,

$$d(\mu, \nu) = \sum_{i=1}^{\infty} \frac{1}{2^i} \left| \int f_i d\mu - \int f_i d\nu \right|.$$

Hence, $\mathcal{P}(X)$ is sequentially compact and μ_n has a convergent subsequence μ_{n_k} , converging to say, μ . We show that μ is σ -invariant: Note $\mu_{n_k} \xrightarrow{w^*} \mu \iff \int f d\mu_{n_k} \rightarrow \int f d\mu$ for all $f \in C(X)$ and that characteristic functions are continuous on X .

$$\begin{aligned} \int f d(\mu \circ \sigma^{-1}) &= \int f \circ \sigma d\mu = \int \lim_{k \rightarrow \infty} f \circ \sigma d\mu_{n_k} \\ &= \int \lim_{k \rightarrow \infty} \frac{1}{2n_k + 1} \sum_{i=-n_k}^{n_k} f(\sigma^{i+1}(x)) \\ &= \lim_{k \rightarrow \infty} \frac{1}{2n_k + 1} \sum_{i=-n_k}^{n_k} f(\sigma^i(x)) \\ &= \int f d\mu. \end{aligned}$$

Hence, $\mu \circ \sigma^{-1} = \mu$. Let $A = \{y \in X \mid y_0 = 1\}$. Then, note $\delta_{\sigma^k(x)}(A) = 1 \iff \sigma^k(x) \in A \iff x_{-k} = 1 \iff -k \in S$. So,

$$\mu(A) = \limsup_{k \rightarrow \infty} \mu_{n_k}(A) = \limsup_{k \rightarrow \infty} \frac{|S \cap \{n_1, \dots, n_k\}|}{2n_k + 1} = d(A) > 0.$$

Applying multiple recurrence theorem, we get for all $k \in \mathbb{N}$, there exists $n \in \mathbb{N}$ such that

$$\mu(A \cap \sigma^{-n}A \cap \sigma^{-2n}A \cap \dots \sigma^{-kn}A) > 0.$$

Since $\mu = \lim \mu_{n_k}$, there exists some $p \in \mathbb{N}$ such that $\mu_p(A \cap \sigma^{-n}A \cap \sigma^{-2n}A \cap \dots \sigma^{-kn}A) > 0$, that is, for some $m \in [-p, p]$,

$$\begin{aligned} \delta_{\sigma_{n^k}(x)}(A \cap \sigma^{-n}A \cap \sigma^{-2n}A \cap \dots \sigma^{-kn}A) &= 1 \\ \implies x &\in \bigcap_{i=1}^k \sigma^{-m-ni}(A) \\ \implies \{\sigma^{m+ni}(x)\}_{i=1}^k &\in A \\ \implies \{-m-ni\}_{i=1}^k &\in S. \end{aligned}$$

□

3.2 Outline of proof of Multiple Recurrence

The rest of the chapter will deal with understanding the proof of the multiple recurrence theorem. We prove a stronger statement, and call this property the *MR-property*:

Theorem 3.2.1. *For a measure-preserving system (X, \mathcal{B}, μ, T) , and a set $A \in \mathcal{B}$, with $\mu(A) > 0$, and for any $k \in \mathbb{N}$, we have*

$$\liminf_{N \rightarrow \infty} \frac{1}{N} \sum_{n=1}^N \mu(A \cap T^{-n}A \cap T^{-2n}A \cap \dots T^{-kn}A) > 0$$

We first show that MR-property is satisfied by some systems, such as, weak-mixing systems and Kronecker systems. The main ingredient is understanding how the property lifts by weak-mixing and compact extensions, and understanding the relationship between the two extensions. We then lift the MR-property from Kronecker systems to arbitrary measure-preserving systems, using appropriate extensions.

Definition 3.2.2 (Extensions and factors). Let $(X, \mathcal{B}_X, \mu, T)$ and $(Y, \mathcal{B}_Y, \nu, S)$ be two MPS. Y is a *factor* of X if there are sets $X' \in \mathcal{B}_X$, $Y' \in \mathcal{B}_Y$ with $\mu(X') = 1$, $\nu(Y') = 1$, $TX' \subset X'$, $SY' \subset Y'$, and a measure-preserving map $\phi : X' \rightarrow Y'$, such that $\phi \circ T = S \circ \phi$. And X is called an *extension* of Y .

Essentially, we have a measure-preserving map $\phi : X \rightarrow Y$ such that $\phi \circ T = S \circ \phi$ a.e. We next show that factors of a system are in 1-1 correspondence with T -invariant sub- σ -algebras ($\mathcal{A} \subset \mathcal{B}_X$ is T -invariant if $T^{-1}\mathcal{A} = \mathcal{A} \mod \mu$). Given a factor $\phi : X \rightarrow Y$, we have T -invariant sub- σ -algebras, $\phi^{-1}\mathcal{B}_Y$. For the other direction,

Theorem 3.2.3. *If (X, \mathcal{B}, μ, T) is a measure-preserving system and $\mathcal{A} \subset \mathcal{B}_X$ is a T -invariant sub- σ -algebra, then*

Proof. □

We prove the multiple recurrence theorem in the following steps, which we elaborate on in the subsequent sections.

Step 1.

- (i) Every MPS has an invertible extension
- (ii) A MPS has MR-property if and only if its invertible extension does.
- (iii) Every invertible extension has a factor, which is a standard probability space.

Now, given any MPS (X, \mathcal{B}, μ, T) , consider its invertible extension and then a standard probability factor of it. If the standard probability factor satisfies MR-property, then (X, \mathcal{B}, μ, T) also does. Thus, it is sufficient to prove that any standard measure-preserving system (X, \mathcal{B}, μ, T) satisfies MR-property.

Step 2.

- (i) Weak-mixing and compact systems satisfy MR-property.
- (ii) A system is not weak-mixing if and only if it has a non-trivial compact factor.
- (iii) If $\mathcal{A}_1 \subset \mathcal{A}_2 \subset \dots$ is an increasing chain of factors of X that satisfy MR-property, then $\sigma(\cup_{n \geq 1} \mathcal{A}_n)$ also satisfies it.

So, if X is weak-mixing, we are done. If it is not weak-mixing, then consider the family of factors, $\mathcal{A} \subset \mathcal{B}$ such that (X, \mathcal{A}, μ, T) satisfy MR-property. We know, it is non-empty, since there is a compact factor that satisfies it. By Zorn's lemma, we can show there is a maximal sub- σ -algebra of this family, call it, \mathcal{B}_∞ .

Step 3

- (i) If X is a weak-mixing extension of Y , which satisfies MR-property, then X also satisfies MR-property.
- (ii) If X is a compact extension of Y , which satisfies MR-property, then X also satisfies MR-property.
- (iii) If $X \rightarrow Y$ is not a weak-mixing extension, then there exists an intermediate factor of $X \rightarrow Z$, such that $Z \rightarrow Y$ is a compact extension.

Hence, if $(X, \mathcal{B}, \mu, T) \rightarrow (X, \mathcal{B}_\infty, \mu, T)$ is a weak-mixing extension, we are done. And if it is not, there is a non-trivial compact extension $(X, \mathcal{C}) \rightarrow (X, \mathcal{B}_\infty)$. So (X, \mathcal{C}) satisfies MR-property, which contradicts the maximality of \mathcal{B}_∞ . Thus, the extension must be weak-mixing, and this completes the proof.

3.2.1 Step 1

Definition 3.2.4 (Invertible extension). Any (X, \mathcal{B}, μ, T) has an invertible extension, $(\tilde{X}, \tilde{\mathcal{B}}, \tilde{\mu}, \tilde{T})$, where

- $\tilde{X} = \{x \in X^{\mathbb{Z}} \mid Tx_k = x_{k+1} \ \forall k \in \mathbb{Z}\}$
- $(\tilde{T}x)_k = x_{k+1}$ for all $k \in \mathbb{Z}$ and $x \in \tilde{X}$
- $\tilde{\mathcal{B}}$ is the product σ -algebra
- $\tilde{\mu}$ is the product measure

$\pi_0 : \tilde{X} \rightarrow X$ is called the invertible extension, where $\pi_0 : X^{\mathbb{Z}} \rightarrow X$ is the 0-th projection.

It is easy to show that the above is a MPS. Many properties like ergodicity, weak-mixing and MR-property are preserved under invertible extensions.

Theorem 3.2.5. (X, \mathcal{B}, μ, T) has MR-property if and only if its invertible extension does.

Proof. Note, $\mu = \tilde{\mu} \circ \pi_0^{-1}$. Let $\tilde{A} \in \tilde{\mathcal{B}}$ with positive measure, and let $A = \pi_0(\tilde{A})$. Then $\pi_0^{-1}(A \cap \dots T^{-kn}A) = (\tilde{A} \cap \dots T^{-kn}\tilde{A})$. This follows from $x \in \pi_0^{-1}(A \cap \dots T^{-kn}A) \iff T^{in}x_0 = x_{in} \in A$, for $0 \leq i \leq k$ and $x_{in} \in A \iff \tilde{T}^{in}x \in \tilde{A}$. Now, the result is clear for one direction. For the converse, start with $A \in \mathcal{B}$ with positive measure and let $\tilde{A} = \pi_0^{-1}(A)$. \square

Theorem 3.2.6. Every invertible system (X, \mathcal{B}, μ, T) has a factor, that is a standard probability space.

Proof. Consider the system, $\{0, 1\}^{\mathbb{Z}}$, with the product σ -algebra and measure, and the shift operator. Fix $A \in \mathcal{B}$ of positive measure and define a map $\phi : X \rightarrow \{0, 1\}^{\mathbb{Z}}$, by $\phi(x) = \chi_A(T^n x)$. ϕ is clearly a factor. \square

3.2.2 Step 2

To define compact systems, we recall some definitions from point-set topology. A set is *precompact* if its closure is compact and is *totally bounded* if for all $\varepsilon > 0$, the set can be covered by a finite number of ε -balls. The two are equivalent in a complete metric space.

Definition 3.2.7 (Almost periodic functions). $f \in L^2(X)$ is *almost periodic* if its orbit $\{T^n f\}_{n \in \mathbb{Z}}$ is precompact in $L^2(X)$.

Definition 3.2.8 (Compact or almost periodic systems). (X, \mathcal{B}, μ, T) is a *compact system* or *almost periodic* if every $f \in L^2(X)$ is almost periodic.

Theorem 3.2.9. *Compact systems satisfy MR-property.*

Proof. □

Theorem 3.2.10. (X, \mathcal{B}, μ, T) is not weak-mixing if and only if it has a non-trivial compact factor.

Proof. □

Theorem 3.2.11. Let (X, \mathcal{B}, μ, T) be a standard invertible MPS and $\mathcal{A}_1 \subset \mathcal{A}_2 \subset \dots$ is an increasing chain of factors that satisfy MR-property, then $\sigma(\cup_{n \geq 1} \mathcal{A}_n)$ also satisfies it.

Proof. Let $\mathcal{A} = \sigma(\cup A_i)$ and $A \in \mathcal{A}$ with positive measure. Then for any $\varepsilon > 0$, there exists some $A_1 \in \mathcal{A}_n$ for some $n \in \mathbb{N}$ such that $\mu(A \triangle A_1) < \varepsilon$. Fix $k \in \mathbb{N}$ and let $\eta = \frac{1}{2(k+1)}$ and $\varepsilon = \frac{1}{4}\eta\mu(A)$. Define $A_0 = \{y \in A_1 \mid \mu_x(A) \geq 1 - \eta\}$.

Claim: $\mu(A_0) > \frac{1}{2}\mu(A)$.

Proof.

$$\frac{1}{4}\eta\mu(A) = \varepsilon > \mu(A_1 \setminus A) = \int_{A_1} \mu_x(A_1 \setminus A) d\mu \geq \int_{A_1 \setminus A_0} (1 - \mu_x(A)) d\mu \geq \eta\mu(A_1 \setminus A_0).$$

Hence, $\mu(A_1 \setminus A_0) \leq \frac{1}{4}\mu(A)$, which implies $\mu(A_0) = \mu(A_1) - \mu(A_1 \setminus A_0) > \frac{1}{2}\mu(A)$ (we can show $\mu(A_1) \geq \frac{3}{4}\mu(A)$). □

Claim: $\mu(A \cap \dots T^{-kn} A) \geq \frac{1}{2}\mu(A_0 \cap \dots T^{-kn} A_0)$.

Proof. If $x \in A_0 \cap \dots T^{-kn} A_0$, then for $0 \leq i \leq k$,

$$\mu_x(T^{-in}(A)) = \mu_{T^{in}x}(A) \geq 1 - \eta.$$

Thus,

$$\mu_x(A \cap \dots T^{-kn} A) \geq 1 - (k+1)\eta = \frac{1}{2}.$$

The claim follows by integrating over $A_0 \cap \dots T^{-kn} A_0$. □

Finally, since A_0 satisfies MR-property, by the above claim, A also satisfies. □

3.2.3 Weak-mixing extension

Definition 3.2.12 (Weak-mixing extension). Let $\phi : (X, \mathcal{B}_X, \mu, T) \rightarrow (Y, \mathcal{B}_Y, \nu, S)$ be an extension. Define a measure $\mu \times_Y \mu$ on $(X \times X, \mathcal{B} \times \mathcal{B})$ given by disintegration with respect to $\mathcal{A} := \phi^{-1}(\mathcal{B}_Y) \subset \mathcal{B}_X$,

$$(\mu \times_Y \mu)_x = \mu_x \times \mu_x.$$

Then $(\tilde{X}, \tilde{\mathcal{B}}, \tilde{\mu}, \tilde{T}) := (X \times X, \mathcal{B} \times \mathcal{B}, \mu \times_Y \mu, T \times T)$ is a MPS and the extension is called *weak-mixing* if this system is ergodic.

Note, if we take the trivial system $(X, \{\phi, X\}, \mu, T)$, then any weak-mixing extension over it is simply a weak-mixing system. This follows from the fact that the measure disintegration over the trivial σ -algebra is the measure itself, i.e., $\mu_x = \mu$.

Theorem 3.2.13. *If X is a weak-mixing extension of Y , which satisfies MR-property, then X also satisfies MR-property.*

Corollary 3.2.14. *If X is weak-mixing, then it satisfies MR-property.*

Proof. Follows from the remark above. \square

We fix some notations for this section. We have $\phi : (X, \mathcal{B}_X, \mu, T) \rightarrow (Y, \mathcal{B}_Y, \nu, S)$ a weak-mixing extension. Let $\mathcal{A} := \phi^{-1}(\mathcal{B}_Y)$. For $f, g \in L^\infty(X)$, define $f \otimes g \in L^\infty(X \times X)$,

$$f \otimes g(x_1, x_2) = f(x_1)g(x_2).$$

We start with a useful identity.

Lemma 3.2.15. *If $f, g \in L^\infty_\mu(X)$, then*

$$\int f \otimes g d\tilde{\mu} = \int \mathbb{E}(f|\mathcal{A}) \mathbb{E}(g|\mathcal{A}) d\mu. \quad (3.1)$$

Proof.

$$\begin{aligned} \int f \otimes g(x_1, x_2) d\tilde{\mu}(x_1, x_2) &= \iint f \otimes g(x_1, x_2) d\mu_y(x_1, x_2) d\mu(y) \\ &= \iint \int f(x_1)g(x_2) d\mu_y(x_1) d\mu_y(x_2) d\mu(y) \\ &= \int \mathbb{E}(f|\mathcal{A})(y) \mathbb{E}(g|\mathcal{A})(y) d\mu(y). \end{aligned}$$

\square

Lemma 3.2.16. *If $f, g \in L^\infty_\mu(X)$, then*

$$\lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \|\mathbb{E}(fT^n g|\mathcal{A}) - \mathbb{E}(f|\mathcal{A}) \mathbb{E}(T^n g|\mathcal{A})\|_{L^2_\mu} = 0.$$

Proof. First let f be such that $\mathbb{E}(f|\mathcal{A}) = 0$, then

$$\begin{aligned} \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int \mathbb{E}(fT^n g|\mathcal{A})^2 d\mu &= \lim_{n \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int (f \otimes f)(T^n g \otimes T^n g) d\tilde{\mu} \\ &= \int (f \otimes f) \left[\lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \tilde{T}^n(g \otimes g) \right] d\tilde{\mu} \\ &= \int (f \otimes f) d\tilde{\mu} \int (g \otimes g) d\tilde{\mu} = 0. \end{aligned}$$

The last line is using Birkhoff's ergodic theorem (by ergodicity of \tilde{T}). Now, for any $f \in L^\infty(X)$, consider $f - \mathbb{E}(f|\mathcal{A})$. Then $\mathbb{E}(f - \mathbb{E}(f|\mathcal{A})|\mathcal{A}) = 0$, and

$$\mathbb{E}((f - \mathbb{E}(f|\mathcal{A})) \cdot T^n g|\mathcal{A}) = \mathbb{E}(fT^n g|\mathcal{A}) - \mathbb{E}(f|\mathcal{A}) \mathbb{E}(T^n g|\mathcal{A}).$$

\square

Lemma 3.2.17 (van der Corput lemma). *If $\{x_n\}$ be a bounded sequence in a Hilbert space H , and*

$$d\text{-}\lim_{h \rightarrow \infty} \left(\lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \langle x_n, x_{n+h} \rangle \right) = 0,$$

then $\lim_{N \rightarrow \infty} \left\| \frac{1}{N} \sum_{n=0}^{N-1} x_n \right\| = 0$.

Proof. We normalize so that $\|x_n\| = 1$, which we can write as $x_n = O(1)$. This gives,

$$\frac{1}{N} \sum_{n=0}^{N-1} x_n = \frac{1}{N} \sum_{n=0}^{N-1} x_{n+h} + O\left(\frac{|h|}{N}\right).$$

Averaging over $0 \leq h \leq H-1$ for some $H \geq 1$,

$$\frac{1}{N} \sum_{n=0}^{N-1} x_n = \frac{1}{N} \sum_{n=0}^{N-1} \frac{1}{H} \sum_{h=0}^{H-1} x_{n+h} + O\left(\frac{H}{N}\right).$$

By triangle inequality and squaring (using $(a+b)^2 \leq 2(a^2+b^2)$),

$$\begin{aligned} \left\| \frac{1}{N} \sum_{n=0}^{N-1} x_n \right\|^2 &\leq O\left(\frac{1}{N} \sum_{n=0}^{N-1} \left\| \frac{1}{H} \sum_{h=0}^{H-1} x_{n+h} \right\|^2 \right) + O\left(\frac{H^2}{N^2}\right) \\ &= O\left(\frac{1}{H^2} \sum_{0 \leq h, h' \leq H} \frac{1}{N} \sum_{n=0}^{N-1} \langle x_{n+h}, x_{n+h'} \rangle \right) + O\left(\frac{H^2}{N^2}\right). \end{aligned}$$

Keeping H fixed,

$$\limsup_{N \rightarrow \infty} \left\| \frac{1}{N} \sum_{n=0}^{N-1} x_n \right\|^2 \leq O\left(\frac{1}{H^2} \sum_{0 \leq h, h' \leq H} \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \langle x_{n+h}, x_{n+h'} \rangle \right).$$

Rewriting, $\sum_{n=0}^{N-1} \langle x_{n+h}, x_{n+h'} \rangle = \sum_{n=0}^{N-1} \langle x_{n+|h-h'|}, x_n \rangle$ gives

$$\limsup_{N \rightarrow \infty} \left\| \frac{1}{N} \sum_{n=0}^{N-1} x_n \right\|^2 \leq O\left(\frac{1}{H} \sum_{h=0}^H \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \langle x_{n+h}, x_n \rangle \right).$$

The result follows by taking the limit as $H \rightarrow \infty$. □

Lemma 3.2.18. *For $f_1, \dots, f_k \in L^\infty(X)$,*

$$\lim_{N \rightarrow \infty} \left\| \frac{1}{N} \sum_{n=0}^{N-1} (T^n f_1 \cdots T^{kn} f_k - T^n \mathbb{E}(f_1 | \mathcal{A}) \cdots T^{kn} \mathbb{E}(f_k | \mathcal{A})) \right\|_2 = 0.$$

Proof. Proof is by induction on k . When $k = 1$, the result follows from the mean ergodic theorem: $\frac{1}{N} \sum_{n=0}^{N-1} T^n f \rightarrow \int f$ and $\frac{1}{N} \sum_{n=0}^{N-1} T^n \mathbb{E}(f | \mathcal{A}) \rightarrow \int \mathbb{E}(f | \mathcal{A}) = \int f$ in $\|\cdot\|_2$.

For the inductive step, we can assume $\mathbb{E}(f_j | \mathcal{A}) = 0$ for some j . Using

$$\prod_{i=1}^k a_i - \prod_{i=1}^k b_i = \sum_{j=1}^k \left[\left(\prod_{i=1}^j a_i \right) (a_j - b_j) \left(\prod_{i=j+1}^k b_i \right) \right],$$

and with $a_i = T^{in} f_i$ and $b_i = T^{in} \mathbb{E}(f_i | \mathcal{A})$, we can reduce the original term to a sum with each summand having the same form, but with a function $T^{in} f_i - T^{in} \mathbb{E}(f_i | \mathcal{A})$ that satisfies our assumption.

The result follows if we can use van der Corput lemma, with $x_n = \prod_{i=1}^k T^{in} f_i$.

$$\begin{aligned}
 \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \langle x_n, x_{n+h} \rangle &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int \prod_{i=1}^k (T^{in} f_i \cdot T^{i(n+h)} f_i) \\
 &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int \prod_{i=1}^k T^{in} (f_i \cdot T^{ih} f_i) \\
 (T\text{-invariance of } \mu) &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int f_1 T^h f_1 \prod_{i=2}^k T^{(i-1)n} (f_i \cdot T^{ih} f_i) \\
 (\text{induction hypothesis}) &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int f_1 T^h f_1 \prod_{i=2}^k \mathbb{E} \left(T^{(i-1)n} (f_i \cdot T^{ih} f_i) \middle| \mathcal{A} \right) \\
 &\leq \prod_{i \neq j} \|f_i\|_\infty \|\mathbb{E}(f_j T^{jh} f_j | \mathcal{A})\|_2
 \end{aligned}$$

By the previous lemma, the last term converges in density to 0. This satisfies the hypothesis for van der Corput lemma. \square

Lemma 3.2.19. \tilde{X} is a weak-mixing extension of Y .

Proof. We need to show $(\tilde{X} \times \tilde{X}, \mathcal{B}_{\tilde{X}} \times \mathcal{B}_{\tilde{X}}, \tilde{\mu} \times_Y \tilde{\mu}, \tilde{T} \times \tilde{T})$ is ergodic. We work on a dense set of $L^\infty(\tilde{X} \times \tilde{X})$. Let $F = f_1 \otimes f_2$ and $G = g_1 \otimes g_2$ for $f_i, g_i \in L^\infty_{\tilde{\mu}}(\tilde{X})$. Here, $\mathcal{A} :=$

$$\begin{aligned}
 \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int F(\tilde{T} \times \tilde{T})^n G d(\tilde{\mu} \times_Y \tilde{\mu}) &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int (f_1 T^n g_1) \otimes (f_2 T^n g_2) d(\tilde{\mu} \times_Y \tilde{\mu}) \\
 (\text{using 3.1}) &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int \mathbb{E}(f_1 T^n g_1 | \mathcal{A}) \mathbb{E}(f_2 T^n g_2 | \mathcal{A}) d\mu \\
 (\text{by previous lemma}) &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int \prod_{i=1}^2 \mathbb{E}(f_i | \mathcal{A}) T^n \mathbb{E}(g_i | \mathcal{A}) d\mu \\
 (\text{by ergodicity of } \tilde{T}) &= \int \prod_{i=1}^2 \mathbb{E}(f_i | \mathcal{A}) d\mu \int \prod_{i=1}^2 T^n \mathbb{E}(T^n g_i | \mathcal{A}) d\mu \\
 &= \int F d\tilde{\mu} \int G d\tilde{\mu}.
 \end{aligned}$$

\square

Lemma 3.2.20. For $f_0, \dots, f_k \in L^\infty(X)$,

$$\lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \left\| \mathbb{E} \left(\prod_{i=0}^k T^{in} f_i \middle| \mathcal{A} \right) - \prod_{i=0}^k \mathbb{E}(T^{in} f_i | \mathcal{A}) \right\|_{L^2_\mu} = 0.$$

Proof. We prove by induction on k . First, assume $\mathbb{E}(f_k | \mathcal{A}) = 0$. By the previous lemma, \tilde{X} is a weak-mixing extension over Y . Using lemma 17,

$$\lim_{N \rightarrow \infty} \left\| \left(\frac{1}{N} \sum_{n=0}^{N-1} \prod_{i=1}^k \tilde{T}^{in} (f_i \otimes f_i) \right) \right\|_{L^2_\mu} = 0,$$

from the fact that $\mathbb{E}(f_k \otimes f_k | \mathcal{A}) = 0$ (by application of 3.1). Note,

$$\prod_{i=1}^k \tilde{T}^{in}(f_i \otimes f_i) = \prod_{i=1}^k T^{in} f_i \otimes \prod_{i=1}^k T^{in} f_i,$$

which gives, using 3.1 and Cauchy-Schwarz inequality,

$$\int \left(\frac{1}{N} \sum_{n=0}^{N-1} \mathbb{E}(T^{in} f_i | \mathcal{A}) \right) d\mu = \int \left(\frac{1}{N} \sum_{n=0}^{N-1} \prod_{i=1}^k \tilde{T}^{in}(f_i \otimes f_i) \right) d\tilde{\mu} \rightarrow 0.$$

Now, for any $f_k \in L^\infty(X)$, consider instead $(f_k - \mathbb{E}(f_k | \mathcal{A})) + \mathbb{E}(f_k | \mathcal{A})$. The proof for $f_k - \mathbb{E}(f_k | \mathcal{A})$ is as above. For $\mathbb{E}(f_k | \mathcal{A})$,

$$\begin{aligned} & \frac{1}{N} \sum_{n=0}^{N-1} \left\| \mathbb{E} \left(\prod_{i=0}^{k-1} T^{in} f_i \cdot T^{kn} \mathbb{E}(f_k | \mathcal{A}) \middle| \mathcal{A} \right) d\mu - \prod_{i=0}^{k-1} \mathbb{E}(T^{in} f_i | \mathcal{A}) \cdot T^{kn} \mathbb{E}(f_k | \mathcal{A}) \right\|_2^2 \\ &= \frac{1}{N} \sum_{n=0}^{N-1} \|T^{kn} \mathbb{E}(f_k | \mathcal{A})\|_\infty^2 \left\| \mathbb{E} \left(\prod_{i=0}^{k-1} T^{in} f_i \middle| \mathcal{A} \right) - \prod_{i=0}^{k-1} \mathbb{E}(T^{in} f_i | \mathcal{A}) \right\|_2^2 \rightarrow 0, \end{aligned}$$

using the induction hypothesis. □

Corollary 3.2.21. For $B_0, \dots, B_k \in \mathcal{B}_X$,

$$\frac{1}{N} \sum_{n=0}^{N-1} \int [\mu_x(B_0 \cap \dots \cap B_k) - \mu_x(B_0) \dots \mu_x(B_k)]^2 d\mu \rightarrow 0.$$

We are now ready for the main result.

Theorem 3.2.22. If Y satisfies MR-property, then so does X .

Proof. Let $B \in \mathcal{B}_X$ with positive measure. Choose $a > 0$ such that the set $A = \{x \mid \mu_x(B) \geq a\}$ has positive measure. Then,

$$\begin{aligned} \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \mu(B \cap \dots \cap T^{-kn} B) &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int \mu_x(B \cap \dots \cap T^{-kn} B) d\mu \\ &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int \mu_x(B) \dots \mu_x(T^{-kn} B) d\mu \\ &= \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} \int \mu_x(B) \dots \mu_{T^{kn} x}(B) d\mu \\ &\geq \lim_{N \rightarrow \infty} \frac{1}{N} \sum_{n=0}^{N-1} a^{k+1} \mu(A \cap \dots \cap T^{-kn} A) > 0. \end{aligned}$$

The penultimate inequality comes from restricting the integral to $A \cap \dots \cap T^{-kn} A$, on which we have $\mu_{T^{in} x}(B) > 0$. □

3.2.4 Compact extension

Definition 3.2.23 (Relative almost periodic functions). Let $(X, \mathcal{B}_X, \mu, T) \rightarrow (Y, \mathcal{B}_Y, \nu, S)$ be an extension. $f \in L^2_\mu(X)$ is *almost periodic* relative to Y if for every $\varepsilon > 0$, there is an $r \in \mathbb{Z}$ and function g_1, \dots, g_r such that

$$\min_{i=1, \dots, r} \|T^k f - g_i\|_{L^2_{\mu_x}} < \varepsilon.$$

for all $n \in \mathbb{N}$ and a.e $x \in X$.

Definition 3.2.24 (Compact extension). $X \rightarrow Y$ is a *compact extension*, if the set of functions almost periodic relative to Y is dense in $L^2_\mu(X)$, where μ_x is the measure disintegration .

Theorem 3.2.25. *If X is a compact extension of Y , which satisfies MR-property, then X also satisfies MR-property.*

3.2.5 Step 3.

This is the final theorem to complete the proof.

Theorem 3.2.26. *If $X \rightarrow Y$ is not a weak-mixing extension, then there exists an intermediate factor of $X \rightarrow Z$, such that $Z \rightarrow Y$ is a compact extension.*

Chapter 4

Operator Algebras

This chapter is a brief exposition on the theory of operator algebras.

4.1 C*-Algebras

4.1.1 Definitions

Definition 4.1.1 (Algebra). An *algebra over a field K* is a vector space with an additional multiplication operation, which satisfies the following, for all $x, y, z \in A$ and $a \in \mathbb{C}$,

- (i) distributive: $(x + y)z = xz + yz$ and $z(x + y) = zx + zy$.
- (ii) associative: $(xy)z = x(yz)$.

Definition 4.1.2 (Banach algebra). A *Banach algebra* A is an algebra over \mathbb{C} with a norm with respect to which it is a normed space and is sub-multiplicative: for all $x, y \in A$, $\|xy\| \leq \|x\| \|y\|$.

Definition 4.1.3 (Involution). An *involution* on algebra A is a map $A \ni x \mapsto x^* \in A$ such that

$$(x + y)^* = x^* + y^*, \quad (\lambda x)^* = \bar{\lambda}x^*, \quad (xy)^* = y^*x^*, \quad x^{**} = x.$$

x^* is called *adjoint* of x .

Definition 4.1.4 (C*-algebra). A *C*-algebra* is a Banach algebra with an involution that satisfies the C*-condition:

$$\|x^*x\| = \|x\|^2.$$

Example. 1. $\mathcal{B}(H)$ on any Hilbert space H is a C*-algebra.

- 2. A norm-closed subalgebra of a C*-algebra closed under adjoint is also a C*-algebra. C*-subalgebras of $\mathcal{B}(H)$ are called *concrete C*-algebras*, while the ones in 4.1.4. are called *abstract C*-algebras*.
- 3. $L^\infty(\mathbb{R})$ with pointwise operations and involution ($f \mapsto \bar{f}$).
- 4. $L^1(\mathbb{R})$ with convolution and above involution is not a C*-algebra.
- 5. $C_0(X)$ on locally compact Hausdorff space X . By Gelfand-Naimark theorem, every commutative C*-algebra is isometrically isomorphic to some $C_0(X)$.

4.1.2 Spectrum and positive operators

Definition 4.1.5. Let $a \in \mathcal{B}(\mathcal{H})$,

Resolvent: $\rho(a) = \{\lambda \in \mathbb{C} \mid \lambda - a \text{ is invertible}\}.$

Spectrum: $\sigma(a) = \mathbb{C} \setminus \rho(a).$

Definition 4.1.6 (Positive elements). $a \in \mathcal{B}(\mathcal{H})$ is *positive* if $a^* = a$ and $\sigma(a) \geq 0$.

Theorem 4.1.7. *The following are equivalent:*

- (i) a is positive.
- (ii) $a = b^2$ for some $b \in \mathcal{B}(\mathcal{H})$.
- (iii) $a = x^*x$ for some $x \in \mathcal{B}(\mathcal{H})$.
- (iv) $\langle x\xi, \xi \rangle \geq 0$ for all $\xi \in \mathcal{H}$.

Theorem 4.1.8. *Positive elements span $\mathcal{B}(\mathcal{H})$.*

4.1.3 Representations of C*-algebras

Definition 4.1.9. Let A be a C*-algebra.

Positive functional: A functional $\phi : A \rightarrow \mathbb{C}$ if $\phi(a) > 0$ when $a > 0$.

States: A positive linear functional of norm 1.

Example. $0 < \xi \in H$, define $\phi : \mathcal{B}(H) \rightarrow \mathbb{C}$, $\phi(x) = \langle x\xi, \xi \rangle$ is positive. ϕ is a state if $\|\xi\| = 1$.

Theorem 4.1.10. (i) *Positive linear functionals are continuous.*

- (ii) *If e_λ is an approximate identity on A , then $\|f\| = \lim f(e_\lambda)$.*
- (iii) *A continuous linear functional is a state, if for some approximate identity e_λ , $\|f\| = 1 = \lim f(e_\lambda)$.*

We now look at representations of C*-algebras and the GNS construction. The GNS representation shows that every state is of the form in the above example.

Definition 4.1.11 (Representation). *Representation of A on a Hilbert space H is a *-homomorphism (continuous algebra homomorphism) from A to $\mathcal{B}(H)$.*

Definition 4.1.12 (Cyclic vector). $\xi \in H$ is a *cyclic vector* for the representation $\pi : A \rightarrow \mathcal{B}(H)$ if $\{\pi(x)\xi \mid x \in A\}$ is dense in H .

Theorem 4.1.13 (GNS construction). *For any positive linear functional, ϕ on A , there exists a representation π_ϕ on a Hilbert space H_ϕ , with a cyclic vector ξ_ϕ such that $\|\phi\| = \|\xi_\phi\|^2$ and*

$$\phi(x) = \langle \pi_\phi(x)\xi_\phi, \xi_\phi \rangle \quad \forall x \in A.$$

$(H_\phi, \pi_\phi, \xi_\phi)$ is called the GNS triple for ϕ . We denote the vector $\pi_\phi(x)\xi_\phi$ as \hat{x} .

4.2 von Neumann Algebras

4.2.1 Operator Topologies

To define the various topologies on $\mathcal{B}(\mathcal{H})$, we first describe the topology induced by a family of seminorms.

Definition 4.2.1 (Seminorm). A *seminorm* is a function on a vector space X , $p : X \rightarrow [0, \infty)$ which satisfies all properties of a norm, except positive-definiteness, ie, p satisfies, for $x, y \in X$ and $\alpha \in \mathbb{C}$

- i. $p(x) \geq 0$
- ii. $p(\alpha x) = |\alpha|p(x)$
- iii. $p(x + y) \leq p(x) + p(y)$.

Definition 4.2.2 (Weak topology). $\{f_i : X \rightarrow X_i\}_{i \in I}$ is a family of maps, and τ_i is a topology on X_i with subbases S_i . We define a topology τ on X , called the *weak topology induced by $\{f_i\}$* , by defining a subbases, $S = \{f_i^{-1}(V) \mid V \in S_i\}_{i \in I}$.

We can show the definition is independent of the choice of subbases. Some of the properties of weak topology are as follows.

Theorem 4.2.3. (i) τ is the smallest topology such that the f_i 's are continuous.

- (ii) For a topological space Z and a function $g : Z \rightarrow X$, g is continuous if and only if $f_i \circ g$ is continuous for all $i \in I$.
- (iii) A net $\{x_\lambda\}_{\lambda \in \Lambda}$ converges to x in τ if and only if the net $f_i(x_\lambda)_{\lambda \in \Lambda}$ converges to $f_i(x)$ in τ_i for all $i \in I$.

If \mathcal{B} is a Banach space, \mathcal{B}^* its dual and $\mathcal{F} \subset \mathcal{B}^*$ a vector subspace, then we denote the weak topology on \mathcal{B} induced by \mathcal{F} as $\sigma(\mathcal{B}, \mathcal{F})$.

Definition 4.2.4 (Topology induced by seminorms). If X is a vector space, and $\{p_i \mid i \in I\}$ is a family of seminorms on X . For each $x \in X$ and $i \in I$, define linear forms $f_{i,x} : X \rightarrow [0, \infty)$, $f_{i,x}(y) = p_i(y - x)$. Then the topology τ on X induced by these linear forms is called the *topology induced by seminorms*.

The following results follow from the previous theorem.

Theorem 4.2.5. (i) For each $x \in X$, $i \in I$, $\varepsilon > 0$. define $U_{(i,x,\varepsilon)} = \{y \in X \mid p_i(y - x) < \varepsilon\}$. The family $\{U_{(i,x,\varepsilon)} \mid x \in X, i \in I, \varepsilon > 0\}$ forms a subbases for τ .

- (ii) A net $\{x_\lambda\}_{\lambda \in \Lambda}$ converges to x in τ if and only if the net $\{p_i(x_\lambda - x)\}_{\lambda \in \Lambda}$ converges to 0 in \mathbb{R} for all $i \in I$.
- (iii) (X, τ) is a topological vector space.
- (iv) τ is the smallest topology such that (X, τ) is a topological vector space with p_i continuous for all $i \in I$.

We now define some topologies on $\mathcal{B}(\mathcal{H})$.

Definition 4.2.6. *Weak operator or wo-topology:* is the topology induced by the family of seminorms

$$x \mapsto |\langle x\xi, \eta \rangle| \text{ for } \xi, \eta \in \mathcal{H}.$$

Strong operator or so-topology: is the topology induced by the family of semi-norms

$$x \mapsto \|x\xi\| \text{ for } \xi \in \mathcal{H}.$$

Ultraweak or σ -weak or w -topology: is the topology induced by the family of semi-norms

$$x \mapsto \left| \sum_{k=1}^{\infty} \langle x\xi_k, \eta_k \rangle \right| \text{ for } \{\xi_k\}, \{\eta_k\} \in \mathcal{H}, \sum_{k=1}^{\infty} \|\xi_k\|^2 < \infty, \sum_{k=1}^{\infty} \|\eta_k\|^2 < \infty.$$

Ultrastrong or σ -strong: is the topology induced by the family of semi-norms

$$x \mapsto \left(\sum_{k=1}^{\infty} \|x\xi_k\|^2 \right)^{1/2} \text{ for } \{\xi_k\} \in \mathcal{H}, \sum_{k=1}^{\infty} \|\xi_k\|^2 < \infty.$$

For $\xi, \eta \in \mathcal{H}$, define the linear forms on $\mathcal{B}(\mathcal{H})$, $\omega_{\xi, \eta}(x) = \langle x\xi, \eta \rangle$. Let $\mathcal{B}(\mathcal{H})_{\sim}$ be the vector space generated by these forms in $\mathcal{B}(\mathcal{H})^*$ and $\mathcal{B}(\mathcal{H})_*$ be the norm closure of $\mathcal{B}(\mathcal{H})_{\sim}$ in $\mathcal{B}(\mathcal{H})^*$. Then w -topology is also $\sigma(\mathcal{B}(\mathcal{H}), \mathcal{B}(\mathcal{H})_{\sim})$. And we can show w -topology is given by $\sigma(\mathcal{B}(\mathcal{H}), \mathcal{B}(\mathcal{H})_*)$. This follows from the density of finite-rank operators in trace-class operators.

We have the following relations between the topologies, ultraweak \subset weak, ultrastrong \subset strong, weak \subset strong, ultraweak \subset ultrastrong and norm is the strongest topology. We shall discuss some more properties of these topologies.

Lemma 4.2.7. *Let \mathcal{B} be a vector space, ϕ be a linear form, p_1, p_2, \dots, p_n be a seminorm on \mathcal{B} such that, for $x \in \mathcal{B}$, $|\phi(x)| \leq \sum_{k=1}^n p_k(x)$. Then there exists linear forms ϕ_1, \dots, ϕ_n such that*

$$\phi = \sum_{k=1}^n \phi_k \text{ and } |\phi_k(x)| \leq p_k(x), x \in \mathcal{B}, k = 1, 2, \dots, n.$$

Theorem 4.2.8. *Let \mathcal{B} be a Banach space, $\mathcal{F} \subset \mathcal{B}^*$ a vector subspace and ϕ a linear form on \mathcal{B} . With the norm topology on \mathcal{B}^* , denote the closed unit ball of \mathcal{B} by \mathcal{B}_1 and the closure of \mathcal{F} by $\overline{\mathcal{F}}$, then*

- (i) ϕ is $\sigma(\mathcal{B}, \mathcal{F})$ -continuous $\iff \phi \in \mathcal{F}$.
- (ii) ϕ is $\sigma(\mathcal{B}, \mathcal{F})$ -continuous on $\iff \phi \in \mathcal{F}$.
- (iii) $\sigma(\mathcal{B}, \mathcal{F})$ and $\sigma(\mathcal{B}, \overline{\mathcal{F}})$ coincide on \mathcal{B}_1 .

Corollary 4.2.9. *Let \mathcal{H} be a Hilbert space. Then*

- (i) $\mathcal{B}(\mathcal{H})_{\sim}$ is the set of all w -continuous linear forms on $\mathcal{B}(\mathcal{H})$.
- (ii) $\mathcal{B}(\mathcal{H})_*$ is the set of all w -continuous linear forms on $\mathcal{B}(\mathcal{H})$.
- (iii) w -topology and w -topology coincide in $\mathcal{B}(\mathcal{H})_1$.
- (iv) A linear form ϕ on $\mathcal{B}(\mathcal{H})$ is w -continuous \iff its restriction to $\mathcal{B}(\mathcal{H})_1$ is w -continuous.

4.2.2 Operators and projections

The set of all self-adjoint operators is denoted by $\mathcal{B}(\mathcal{H})_{\text{sa}}$ and the set of all projections on $\mathcal{B}(\mathcal{H})$ is denoted by $\mathcal{P}_{\mathcal{B}(\mathcal{H})}$.

On the set of all self-adjoint operators, $\mathcal{B}(\mathcal{H})_{\text{sa}}$, we have an order relation:

$$x, y \in \mathcal{B}(\mathcal{H})_{\text{sa}}, x \geq y \iff x - y \geq 0.$$

Theorem 4.2.10. Let $\{x_i\} \subset \mathcal{B}(\mathcal{H})_{sa}$ be a bounded increasing net. Then, there is an $x \in \mathcal{B}(\mathcal{H})_{sa}$ such that $x = \sup_i x_i$. Moreover, $x = \text{so-lim}_i x_i$.

Lemma 4.2.11. Let $x \in \mathcal{B}(\mathcal{H})$, $0 \leq x \leq 1$, and $e \in \mathcal{P}_{\mathcal{B}(\mathcal{H})}$, then $x \leq e \iff x = xe$.

Definition 4.2.12. Let $\{e_i\}_{i \in I} \subset \mathcal{P}_{\mathcal{B}(\mathcal{H})}$. Define

- $\bigvee_{i \in I} e_i = \text{projection onto } \overline{\sum_{i \in I} e_i \mathcal{H}}$.
It is the least upper bound of the family $\{e_i\}$
- $\bigwedge_{i \in I} e_i = \text{projection onto } \bigcap_{i \in I} e_i \mathcal{H}$.
It is the greatest lower bound of the family $\{e_i\}$

Theorem 4.2.13 (Borel functional calculus). $x \in \mathcal{B}(\mathcal{H})_{sa}$, then we have a $*$ -homomorphism

$$\mathcal{B}(\sigma(x)) \ni f \mapsto f(x) \in \mathcal{R}(\{x\}),$$

where $\mathcal{R}(\{x\})$ is the von Neumann algebra generated by x .

Corollary 4.2.14. A von Neumann algebra equals the norm-closed linear span of its projections.

4.2.3 Elementary von Neumann Algebra

Definition 4.2.15 (von Neumann algebra). A subalgebra $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ is a *von Neumann algebra* if it is unital, self-adjoint and so-closed.

Definition 4.2.16 (Commutant). Let $\mathcal{X} \subset \mathcal{B}(\mathcal{H})$, then commutant of \mathcal{X} , $\mathcal{X}' = \{x' \in \mathcal{B}(\mathcal{H}) \mid x'x = xx' \text{ for all } x \in \mathcal{X}\}$

Theorem 4.2.17 (von Neumann density theorem). Let $\mathcal{A} \subset \mathcal{B}(\mathcal{H})$ be a unital $*$ -subalgebra, then

$$\mathcal{A}'' = \overline{\mathcal{A}}^{wo} = \overline{\mathcal{A}}^{so} = \overline{\mathcal{A}}^{\sigma-wo}$$

If $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ is a von Neumann algebra, we denote the set of projections in \mathcal{M} is denoted by $\mathcal{P}_{\mathcal{M}}$.

Corollary 4.2.18. Let $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ be a von Neumann algebra and $x \in \mathcal{B}(\mathcal{H})$, then the following statements are equivalent,

1. $x \in \mathcal{M}$.
2. $xe' = e'x$ for any projection $e' \in \mathcal{M}'$.
3. $u'^*xu' = x$ for any unitary operator $u' \in \mathcal{M}'$.

Corollary 4.2.19. Let $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ be a von Neumann algebra and $x \in \mathcal{B}(\mathcal{H})$, then $l(x), r(x) \in \mathcal{M}$.

Corollary 4.2.20. Let $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ be a von Neumann algebra, then $\mathcal{P}_{\mathcal{M}}$ is a complete lattice.

Corollary 4.2.21. Let $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ be a von Neumann algebra and $\xi \in \mathcal{H}$, denote

$$p_\xi = [\mathcal{M}'\xi], \text{ and } p'_\xi = [\mathcal{M}\xi].$$

Then, $p_\xi \in \mathcal{M}$, and $p'_\xi \in \mathcal{M}'$.

Definition 4.2.22 (Separating vector). If $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$, then $\xi \in \mathcal{H}$ is a *separating vector* if for $x \in \mathcal{M}$, $x\xi = 0 \implies x = 0$.

Lemma 4.2.23. *If $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$, then $\xi \in \mathcal{H}$ is a separating (respectively cyclic) vector for \mathcal{M} if and only if ξ is a cyclic (respectively separating) vector for \mathcal{M}' .*

Definition 4.2.24 (Central support). Let $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ be a von Neumann algebra and $x \in \mathcal{B}(\mathcal{H})$. The smallest element of the set of projections $p \in \mathcal{M}$ such that $px = x$ is called *central support* of x and denoted by $z(x)$.

Theorem 4.2.25. *Let $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ be a von Neumann algebra and $x \in \mathcal{B}(\mathcal{H})$. Then $z(x) = [(\mathcal{M}x)\mathcal{H}]$.*

Theorem 4.2.26 (Kaplansky density theorem). *Let $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ be a von Neumann algebra. If $\mathcal{A} \subset \mathcal{M}$ be a so-dense $*$ -subalgebra of \mathcal{M} , then the unit ball of \mathcal{A} is so-dense in the unit ball of \mathcal{M} .*

Corollary 4.2.27. *Let $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ be a unital $*$ -algebra. Then \mathcal{M} is a von Neumann algebra if and only if \mathcal{M}_1 is w -compact.*

Corollary 4.2.28. *Let $\phi : \mathcal{M}_1 \rightarrow \mathcal{M}_2$ be a unital w -continuous $*$ -homomorphism between von Neumann algebra, then $\phi(\mathcal{M}_1)$ is a von Neumann algebra.*

Theorem 4.2.29. *Let \mathcal{M} be a von Neumann algebra and \mathcal{N} be a left ideal, then there exists a unique projection $e \in \mathcal{M}$ such that $\overline{\mathcal{N}}^w = \mathcal{M}e$.*

We end with a discussion on support projection of normal forms. This will be used in the next chapter.

Definition 4.2.30 (Support projection of normal form). Let $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ be a von Neumann algebra and ϕ be a normal form on \mathcal{M} .

Theorem 4.2.31. *Let \mathfrak{A} be a C^* -algebra and ϕ be a state on \mathfrak{A} , then the support projection $s(\phi)$ of ϕ in the double dual \mathfrak{A}^{**} is central if and only if $\xi_\phi \in \mathcal{H}_\phi$ is cyclic for $\pi_\phi(\mathfrak{A})'$.*

4.2.4 W^* -algebra

Definition 4.2.32 (W^* -algebra). A C^* -algebra \mathcal{M} is a W^* -algebra if it admits a predual \mathcal{M}_* , that is, \mathcal{M} is isometrically isomorphic to the dual space of some Banach space, which we call the predual of \mathcal{M} and denote by \mathcal{M}_* .

It is true that they are equivalent to von Neumann algebras, which leads us to call W^* -algebras abstract von Neumann algebras. We prove that every von Neumann algebra is a W^* -algebra.

Lemma 4.2.33. *Let X be a normed space and $F \subset X^*$, then the following are equivalent,*

- (i) *The map $\Phi : X \rightarrow F^*$, $\Phi(x)\varphi = \varphi(x)$, where $x \in X, \varphi \in F$ is a linear isometry.*
- (ii) *F separates elements of X and X_1 is $\sigma(X, F)$ -compact.*

Proof. \implies Follows by Banach-Alaoglu theorem.

\impliedby : Since F separates elements of X , Φ is injective. Φ is a contraction:

$$\|\Phi(x)\| = \sup\{|\Phi(x)\phi| \mid \phi \in M_*, \|\phi\| \leq 1\} = \sup\{|\phi(x)| \mid \phi \in M_*, \|\phi\| \leq 1\} \leq \|x\|.$$

Since Φ is w^* -continuous, by the assumption $\Phi(X_1)$ is a w^* -compact convex subset of F^* . Let $f \in (F^*)_1 \setminus \Phi(X_1)$, then by Hahn-Banach separation theorem, there exist $\phi \in F, c, d \in \mathbb{R}$ such that $\operatorname{Re} \phi(x) \leq c < d \leq \operatorname{Re} f(\phi)$ for all $x \in X_1$. The first inequality gives $\|\Phi\| \leq c$, which contradicts the third inequality. Thus, $(F^*)_1 = \Phi(X_1)$, which implies Φ is a surjective isometry. \square

Theorem 4.2.34. *Let $M \subset \mathcal{B}(\mathcal{H})$ be a w -closed subspace and M_* be the space of all w -continuous linear functionals on M . Then for every $\phi \in M_*$ and $\varepsilon > 0$, there exists $\psi \in \mathcal{B}(\mathcal{H})_*$ such that $\phi = \psi|_M$ and $\|\psi\| \leq \|\phi\| + \varepsilon$. Thus, $M_* = \{\psi|_M \mid \psi \in \mathcal{B}(\mathcal{H})_*\}$, $M_* \subset M^*$ is a norm closed subspace and we have a surjective linear isometry $\Phi : M \rightarrow (M_*)^*$, $\Phi(x)\varphi = \varphi(x)$, where $x \in X, \varphi \in M_*$.*

Proof. Let $\psi \in M_*$, then by Hahn-Banach theorem there exists $\theta \in \mathcal{B}(\mathcal{H})_*$ such that $\theta|M = \phi$. Let $K = \{\rho \in \mathcal{B}(\mathcal{H})_* \mid \rho|M = 0\}$ and $d = \inf\{\|\rho - \theta\| \mid \rho \in K\}$. By Hahn-Banach theorem, there is a $f \in \mathcal{B}(\mathcal{H})_*$ such that $\|f\| \leq 1$, $f|K = 0$, $f(\theta) = d$: extend the bounded linear functional on $\mathbb{C}\theta + K$, $\lambda\theta + \rho \mapsto \lambda d$. By the previous lemma, $\mathcal{B}(\mathcal{H})$ is isometric to $(\mathcal{B}(\mathcal{H})_*)^*$. So we have $x \in \mathcal{B}(\mathcal{H})$, $\Phi(x) = f$. Then $\|x\| = \|f\| \leq 1$ and $\Phi(x)\rho = \rho(x) = f(\rho)$, $\rho \in \mathcal{B}(\mathcal{H})_*$. Since $\rho(x) = 0$ for all $\rho \in K$ and M is w-closed, by Hahn-Banach separation theorem, we get $x \in M$. We have

$$d = f(\theta) = \theta(x) = \phi(x) \leq \|\phi\| \|x\| \leq \|\phi\|.$$

Hence, for every $\varepsilon > 0$, there exists $\rho \in K$ such that $\|\theta - \rho\| \leq d + \varepsilon \leq \|\phi\| + \varepsilon$. $\psi = \phi - \rho$ satisfies the first statement. So finally, the wo-topology on M coincides with the weak*-topology and by the previous lemma gives the final statement. \square

This shows that the double dual of C*-algebra is a W*-algebra. This allows us to study some properties of C*-algebras by reducing them to W*-algebras. An important fact to recall the Goldstine's theorem, which states the closed unit ball of a Banach space is w*-dense under the embedding in the closed unit ball of the double dual. So, a Banach space is w*-dense under the embedding in the double dual.

We will now see an instance of reducing C*-algebra problems to W*-algebras.

Definition 4.2.35. Let \mathcal{A} be a C*-algebra, $\mathcal{B} \subset \mathcal{A}$ be a C*-subalgebra, then a linear map $\Phi : \mathcal{A} \rightarrow \mathcal{B}$ is a

Projection: if $\Phi(b) = b$ for every $b \in \mathcal{B}$. Then $\Phi \circ \Phi = \Phi$.

\mathcal{B} -linear: if $\Phi(ab) = \Phi(a)b$ and $\Phi(ba) = b\Phi(a)$ for $a \in \mathcal{A}$, $b \in \mathcal{B}$.

Conditional expectation: if it is \mathcal{B} -linear and a positive map.

Theorem 4.2.36. Every projection of norm 1, $\Phi : \mathcal{A} \rightarrow \mathcal{B}$ is a conditional expectation.

Proof. We first consider the case when \mathcal{A}, \mathcal{B} are W*-algebras. Denote the unit of \mathcal{B} as $1_{\mathcal{B}}$. Let $e \in \mathcal{B}$ be a projection and let $f = 1 - e$. For $x, y \in \mathcal{A}$,

$$\|ex + fy\|^2 = \|(ex + fy)^*(ex + fy)\| = \|x^*ex + y^*fy\| \leq \|x^*ex\| + \|y^*fy\| = \|ex\|^2 + \|fy\|^2.$$

For $\lambda \in \mathbb{R}$,

$$\begin{aligned} (1 + \lambda)^2 \|f\Phi(ex)\|^2 &= \|f\Phi(ex) + \lambda f\Phi(ex)\|^2 = \|f\Phi(ex + \lambda f\Phi(ex))\|^2 \\ &\leq \|ex + \lambda f\Phi(ex)\|^2 \leq \|ex\|^2 + \lambda^2 \|f\Phi(ex)\|^2. \end{aligned}$$

This gives, $0 \leq \|f\Phi(ex)\|^2 \leq \|ex\|^2 - 2\lambda \|f\Phi(ex)\|^2$, which does not hold for arbitrary $\lambda \in \mathbb{R}$ unless $\|f\Phi(ex)\| = 0$. Thus, $(1 - e)\Phi(ex) = f\Phi(ex) = 0$ which implies $\Phi(ex) = e\Phi(ex)$. Interchanging e and f , we get $e\Phi(x - ex) = \Phi(fx) = 0$ which implies $e\Phi(x) = e\Phi(ex)$. Thus,

$$e\Phi(x) = \Phi(ex), \quad x \in \mathcal{A}.$$

Since \mathcal{B} is linearly spanned by its projection, we get $y\Phi(x) = \Phi(yx)$, $y \in \mathcal{B}$. To show $\Phi(xy) = \Phi(x)y$, $y \in \mathcal{B}$, it is sufficient to be able to take adjoint, which we now justify. With $x = 1, e = 1_{\mathcal{B}}$, we have $1_{\mathcal{B}} = \Phi(1_{\mathcal{B}}) = 1_{\mathcal{B}}\Phi(1) = \Phi(1)$. Now, we show Φ is positive. Recall, an element is positive if and only its image is positive under any positive linear functional. Hence, Φ is positive if $\phi = \psi \circ \Phi$ is positive for all positive linear functional ψ on \mathcal{B} . Thus, Φ is self-adjoint.

Finally, we extend to when \mathcal{A}, \mathcal{B} are C*-algebras. Consider the double transpose, $\Phi^{**} : \mathcal{A}^{**} \rightarrow \mathcal{B}^{**}$. By the w-density of \mathcal{A}, \mathcal{B} in $\mathcal{A}^{**}, \mathcal{B}^{**}$ respectively, we can consider \mathcal{B}^{**} to be a W*-subalgebra of \mathcal{A}^{**} and Φ^{**} is also a projection of norm 1 and hence, a conditional expectation. Its restriction to \mathcal{A} will remain a conditional expectation. \square

Definition 4.2.37 (Enveloping von Neumann algebra of a C*-algebra). Let \mathcal{A} be a C*-algebra and consider the universal representation,

$$\pi_{\mathcal{A}} = \bigoplus_{\phi \in S(\mathcal{A})} \pi_{\phi} : \mathcal{A} \rightarrow \mathcal{B}(\mathcal{H}_{\mathcal{A}}).$$

The *enveloping von Neumann algebra* of \mathcal{A} , $N_{\mathcal{A}}$ is defined to be w-closure of $\pi_{\mathcal{A}}(\mathcal{A})$.

Theorem 4.2.38. *There is a surjective linear isometry $\Psi_{\mathcal{A}} : \mathcal{A}^* \rightarrow (N_{\mathcal{A}})_*$ such that $\psi = \Psi_{\mathcal{A}}(\psi) \circ \pi_{\mathcal{A}}$ for $\psi \in \mathcal{A}^*$.*

Proof. Recall every functional on \mathcal{A} can be written as a sum of states uniquely. So, we can define a map on states and extend linearly. If $\phi \in S(\mathcal{A})$, then denote the vector $\zeta_{\phi} \in \mathcal{H}_{\mathcal{A}}$ as the cyclic vector for ϕ , ξ_{ϕ} in the ϕ -th coordinate and 0 otherwise. Define

$$\Psi_{\mathcal{A}} : \mathcal{A}^* \rightarrow (N_{\mathcal{A}})_*, \phi \mapsto \omega_{\zeta_{\phi}}.$$

$\Psi_{\mathcal{A}}$ is an isometry: Let $\sum_{k=1}^n \lambda_k \phi_k \in \mathcal{A}^*$, where $\lambda_k \in \mathbb{C}, \phi_k \in S(\mathcal{A})$. Then $\sum_{k=1}^n \omega_{\zeta_{\phi_k}} \in (N_{\mathcal{A}})_*$ and using Kaplansky density theorem,

$$\begin{aligned} \left\| \sum_{k=1}^n \lambda_k \omega_{\zeta_{\phi_k}} \right\| &= \sup \left\{ \left| \sum_{k=1}^n \lambda_k \omega_{\zeta_{\phi_k}}(\pi_{\mathcal{A}}(x)) \right|, x \in \mathcal{A}_1 \right\} \\ &= \sup \left\{ \left| \sum_{k=1}^n \lambda_k \phi_k(x) \right|, x \in \mathcal{A}_1 \right\} = \left\| \sum_{k=1}^n \lambda_k \phi_k \right\|. \end{aligned}$$

Also, for $\psi \in \mathcal{A}^*$, we have $\Psi_{\mathcal{A}}(\psi)(\pi_{\mathcal{A}}(x)) = \psi(x)$. For surjectivity, let $\theta \in (N_{\mathcal{A}})_*$ and define $\psi_{\mathcal{A}}^*$, $\psi(x) = \theta(\pi_{\mathcal{A}}(x))$, $x \in \mathcal{A}$. By the previous equality and density of $\pi_{\mathcal{A}}(\mathcal{A}) \subset N_{\mathcal{A}}$, we get $\Psi_{\mathcal{A}}(\psi) = \theta$. \square

Theorem 4.2.39. *There is a map $N_{\mathcal{A}} \rightarrow \mathcal{A}^{**}$ which is a surjective linear isometry and a $(w, \sigma(\mathcal{A}^{**}, \mathcal{A}^*))$ -homeomorphism.*

Proof. Using the identification $N_{\mathcal{A}} \cong ((N_{\mathcal{A}})_*)^*$, we find the transpose of the map $\Psi_{\mathcal{A}} : \mathcal{A}^* \rightarrow (N_{\mathcal{A}})_*$ is $\Psi_{\mathcal{A}}^* : N_{\mathcal{A}} \rightarrow \mathcal{A}^{**}$, $\Psi_{\mathcal{A}}^*(x)(\psi) = \psi(\Psi_{\mathcal{A}}(x)) = \Psi_{\mathcal{A}}(x)\psi$. Since the transpose of an isometry is an isometry, $\Psi_{\mathcal{A}}^*$ is also an isometry. It is also a homeomorphism with the w-topology on $N_{\mathcal{A}}$ and weak *-topology on \mathcal{A}^{**} , using the fact that transpose is w*-continuous, transpose of a bijective map is bijective and the inverse is w*-continuous on the unit ball. Hence, we have a $(w, \sigma(\mathcal{A}^{**}, \mathcal{A}^*))$ -homeomorphism. \square

Theorem 4.2.40. *Let \mathcal{M} be a von Neumann algebra with predual \mathcal{M}_* , there exists a unique central projection $p \in N_{\mathcal{M}}$ such that the map $\mathcal{M} \ni x \mapsto \pi_{\mathcal{M}}(x)p \in (N_{\mathcal{M}})p$ is a surjective *-isomorphism.*

Proof. By theorem 4.2.18., we have a linear isometry $\Psi_{\mathcal{M}} : \mathcal{M}^* \rightarrow (N_{\mathcal{M}})^*$ such that $\psi = \Psi_{\mathcal{M}}(\psi) \circ \pi_{\mathcal{M}}$. Then $\Psi|(\mathcal{M}_*)$ is an injective linear isometry from \mathcal{M}_* to $(N_{\mathcal{M}})_*$. Consider its dual $\Phi = (\Psi|(\mathcal{M}_*))^* : N_{\mathcal{M}} = ((N_{\mathcal{M}})_*)^* \rightarrow (\mathcal{M}_*)^* = \mathcal{M}$. Then Φ is a map of norm 1 and hence, $\pi_{\mathcal{M}} \circ \Phi$ is of norm 1. Φ is also a projection from $N_{\mathcal{M}}$ to $\pi_{\mathcal{M}}(\mathcal{M})$: for $x \in \mathcal{M}$ and $\psi \in \mathcal{M}_*$, using the identification $(\mathcal{M}_*)^* = \mathcal{M}$ and some abuse of notation,

$$\Phi(\pi_{\mathcal{M}}(x))(\psi) = \pi_{\mathcal{M}}(x)(\Psi(\psi)) = \Psi(\psi)\pi_{\mathcal{M}}(x) = \psi(x) = x(\psi).$$

This gives $\Phi \circ \pi_{\mathcal{M}} = \text{id}$ and $(\pi_{\mathcal{M}} \circ \Phi)\pi_{\mathcal{M}}(x) = \pi_{\mathcal{M}}(x)$. Thus, $\pi_{\mathcal{M}} \circ \Phi$ is positive and $\pi_{\mathcal{M}}(\mathcal{M})$ -linear. It is easy to show, using $\pi_{\mathcal{M}}(\mathcal{M})$ -linearity and w-density of $\pi_{\mathcal{M}}(\mathcal{M})$ in $N_{\mathcal{M}}$, that $\ker \Phi$ is a 2-sided ideal of $N_{\mathcal{M}}$. Since, dual maps are w-continuous, $\ker \Phi$ is w-closed in $N_{\mathcal{M}}$. Hence, there exists a

unique central projection $p \in N_{\mathcal{M}}$ such that $\ker \Phi = N_{\mathcal{M}}(1 - p)$. And, Φ is a $*$ -homomorphism: for $x, y \in N_{\mathcal{M}}$, using $\Phi \circ \pi_{\mathcal{M}} = \text{id}$, we get $x - (\pi_{\mathcal{M}} \circ \Phi)(x) \in \ker \Phi$, so $xy - (\pi_{\mathcal{M}} \circ \Phi)(x)y \in \ker \Phi$, which gives $\Phi(xy) = \Phi(x)\Phi(y)$. For surjectivity: let $x \in \mathcal{M}$, note $\Phi(\pi_{\mathcal{M}}(x)(1 - p)) = 0 \implies \Phi(\pi_{\mathcal{M}}(x)p) = \Phi(\pi_{\mathcal{M}}(x)) = x$. Thus, $\Phi|_{(N_{\mathcal{M}})p}$ is an isomorphism from $(N_{\mathcal{M}})p$ onto \mathcal{M} . Finally, by w-density of $\pi_{\mathcal{M}}(\mathcal{M})$ in $N_{\mathcal{M}}$ and normality of Φ , we get $\Phi(xp) = x$ for $x \in \mathcal{M}$. \square

Theorem 4.2.41. *Let \mathcal{A} be a C^* -algebra, \mathcal{M} be a von Neumann algebra, and $\Theta : \mathcal{A} \rightarrow \mathcal{M}$ be a bounded linear map, then Θ can be uniquely extended to a normal linear map from $\tilde{\Theta} : \mathcal{A}^{**} \rightarrow \mathcal{M}$.*

Proof. Consider the double dual $\Theta^{**} : \mathcal{A}^{**} \rightarrow \mathcal{M}^{**}$, which is clearly w-continuous. Using the previous two theorems, we get a central projection p and a $*$ -isomorphism $\Phi : \mathcal{M}^{**}p \rightarrow \mathcal{M}$ such that $\Phi(xp) = x$ for $x \in \mathcal{M}$. Define $\tilde{\Theta} : \mathcal{A}^{**} \rightarrow \mathcal{M}$ by $\mathcal{A}^{**} \ni x \mapsto \Phi(\Theta^{**}(x)p) \in \mathcal{M}$. If $x \in \mathcal{M}$, then $\tilde{\Theta}(x) = \Phi(\Theta(x)p) = \Theta(x)$. Hence, $\tilde{\Theta}|_{\mathcal{A}} = \Theta$. And uniqueness follows from w-density of \mathcal{A} in \mathcal{A}^{**} . \square

Chapter 5

Noncommutative Recurrence

5.1 Noncommutative Poincare recurrence

Definition 5.1.1 (C*-dynamical system). A *C*-dynamical system* is a triplet $(\mathfrak{A}, \phi, \alpha)$, where \mathfrak{A} is a C*-algebra, ϕ is a state and $\alpha : \mathfrak{A} \rightarrow \mathfrak{A}$ is a C*-algebra homomorphism.

Definition 5.1.2 (State-preserving C*-dynamical system). A C*-dynamical system $(\mathfrak{A}, \phi, \alpha)$ is *state-preserving* if $\phi \circ \alpha = \phi$.

Example. Given a MPS (X, \mathcal{B}, μ, T) , we have a C*-dynamical system, $(\mathfrak{A}, \phi, \alpha)$, where $\mathfrak{A} = L^\infty_\mu(X)$, $\phi(f) = \int f d\mu$, and $\alpha = U_T$.

We first discuss recurrence in unital C*-algebras to simplify things.

Definition 5.1.3 (Unital C*-dynamical system). A C*-dynamical system is *unital* if \mathfrak{A} is unital, $\phi(1) = 1$ and $\alpha(1) = 1$.

Theorem 5.1.4 (Noncommutative Poincare recurrence theorem). *Let $(\mathfrak{A}, \phi, \alpha)$ be a unital C*-dynamical system. Then*

$$\liminf_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} |\phi(\alpha^k(a^*)a)| > 0,$$

for every $a \in \mathfrak{A}$ with $\phi(a) > 0$.

To prove this, we first prove the noncommutative Khintchine recurrence theorem.

Definition 5.1.5 (Relatively dense set). A subset $N \subset \mathbb{N}$ is *relatively dense* if there is an $L > 0$ such that every interval in \mathbb{N} of length L has an element of N .

Theorem 5.1.6 (Khintchine recurrence theorem). *In (X, \mathcal{B}, μ, T) , for $A \in \mathcal{B}$ and all $\varepsilon > 0$, there exists a relatively dense set $N \subset \mathbb{N}$ such that, for all $n \in N$,*

$$\mu(T^{-n}A \cap A) \geq \mu(A)^2 - \varepsilon.$$

Theorem 5.1.7 (Noncommutative Khintchine recurrence theorem). *In unital C*-dynamical system, $(\mathfrak{A}, \phi, \alpha)$, for every $x \in \mathfrak{A}$ and every $\varepsilon > 0$, there is a relatively dense set $N \subset \mathbb{N}$ such that, for all $n \in N$,*

$$\operatorname{Re} \phi(\alpha^n(x^*)x) \geq |\phi(x)|^2 - \varepsilon.$$

The commutative version follows using the above example, with x as indicator functions. It is also stronger in the sense that it holds for all elements, and not just projections (recall, the projections in $L^\infty(X)$ are the indicator functions).

Lemma 5.1.8. *Let H be a Hilbert space and let $T : H \rightarrow H$ be an operator such that $\|Tx\| = \|x\|$, for all $x \in H$ and $Tv = v$ for some $v \in H$ with norm 1. Then, for every $x \in H$ and $\varepsilon > 0$, there is a relatively dense set $N \subset \mathbb{N}$ such that, for all $n \in n$,*

$$\operatorname{Re} \langle T^n x, x \rangle \geq |\langle x, v \rangle|^2 - \varepsilon.$$

Proof. By the mean ergodic theorem, for all $\varepsilon > 0$, there is an $n \in \mathbb{N}$ such that

$$\left\| \frac{1}{n} \sum_{k=0}^{n-1} T^k x - Px \right\|^2 \leq \frac{\varepsilon}{2},$$

where P is the projection onto the subspace $\{x \in H \mid Tx = x\}$. Denote $x_n = \frac{1}{n} \sum_{k=0}^{n-1} T^k x$. Since $TP = P$, for all $i \in \mathbb{N}$, $\|T^i x_n - Px\|^2 \leq \|T\|^i \|x_n - Px\|^2 \leq \frac{\varepsilon}{2}$. So, for $i \geq n$, $\|T^i x_n - x_n\|^2 \leq 2\varepsilon$, which gives

$$\|x_n\|^2 - \operatorname{Re} \langle T^i x_n, x_n \rangle \leq \varepsilon.$$

Now, $0 \leq \|x_n - \langle x_n, v \rangle v\|^2 = \|x_n\|^2 - |\langle x_n, v \rangle|^2$. By lemma 2.3.1., $Tv = v \implies T^*v = v$. Hence, $\langle x_n, v \rangle = \langle x, v \rangle$. So, $|\langle x, v \rangle|^2 = |\langle x_n, v \rangle|^2 \leq \|x_n\|^2$, and we have, for $i \geq n$,

$$|\langle x, v \rangle|^2 - \varepsilon \leq \operatorname{Re} \langle T^i x_n, x_n \rangle = \frac{1}{n^2} \sum_{j,k=0}^{n-1} \operatorname{Re} \langle T^{i+j} x, T^k x \rangle \leq \frac{1}{n^2} \sum_{j,k=0}^{n-1} \operatorname{Re} \langle T^{i+j-k} x, x \rangle.$$

Thus, for each $m \in \mathbb{N}$, there exists integers $j(m), k(m) \in [0, n-1]$ such that

$$\operatorname{Re} \langle U^{nm+j(m)-k(m)} x, x \rangle \geq |\langle x, v \rangle|^2 - \varepsilon.$$

Take $N := \{nm + j(m) - k(m) \mid m \in \mathbb{N}\}$. It remains to show N is relatively dense. Note, $-n+1 \leq j(m) - k(m) \leq n-1 \implies (m-1)n \leq nm + j(m) - k(m) \leq (m+1)n$. So, each interval of length $2m$ contains an element of N . \square

Definition 5.1.9. Let \mathfrak{A} be a C^* -algebra, ϕ a state and $\alpha : \mathfrak{A} \rightarrow \mathfrak{A}$ a positive linear map such that $\phi \circ \alpha = \phi$. Consider the GNS triple $(H_\phi, \pi_\phi, \xi_\phi)$ for state ϕ . The linear map $\alpha : \mathfrak{A} \rightarrow \mathfrak{A}$ induces a linear map

$$U_\alpha : H_\phi \rightarrow H_\phi, \quad U_\alpha(\widehat{x}) = \widehat{\alpha(x)}.$$

Lemma 5.1.10. *If α is a $*$ -homomorphism, then U_α is an isometry.*

Proof.

$$\|U_\alpha(\widehat{x})\|^2 = \langle \widehat{\alpha(x)}, \widehat{\alpha(x)} \rangle = \phi(\alpha(x)^* \alpha(x)) = \phi(x^* x) = \langle \widehat{x}, \widehat{x} \rangle = \|\widehat{x}\|^2.$$

Thus, U_α is an isometry on $\{\pi_\alpha(x)\xi_\phi\}$, and by density, U_α is well-defined and an isometry on entire H_ϕ . \square

Proof of Theorem 5.1.5. Note, U_α is an isometry and $U_\alpha(\pi_\phi(1)) = \widehat{1}$. We apply the lemma on $H = H_\phi$, $U = U_\alpha$ and $v = \widehat{1}$. Then

$$\operatorname{Re} \langle U_\alpha^n \widehat{x}, \widehat{x} \rangle = \operatorname{Re} \phi(\alpha^n(x^*)x) \geq |\langle \pi_\phi(x), \pi_\phi(1) \rangle|^2 - \varepsilon \geq |\phi(x)|^2 - \varepsilon.$$

\square

Proof of Theorem 5.1.2. We use lemma 2.4.7. (i) \iff (iii) on theorem 5.1.5. We need to show the set N has positive density: if every interval of length L has an element of N , $|N \cap [0, nL]| \geq n$. So, $d(N) \geq 1/L > 0$. \square

We now consider non-unital C^* -algebras.

Lemma 5.1.11. *Let \mathfrak{A} be a C^* -algebra, ϕ a state and $\alpha : \mathfrak{A} \rightarrow \mathfrak{A}$ a positive linear map such that $\phi \circ \alpha = \phi$.*

1. *If $\phi(\alpha(x)^*\alpha(x)) \leq \phi(x^*x)$ for all $x \in \mathfrak{A}$, then*

(i) U_α is a contraction.

(ii) $U_\alpha \xi_\phi = \xi_\phi$.

(iii) *If P is the orthogonal projection onto $\{\xi \in H_\phi \mid U_\alpha \xi = \xi\}$, then*

$$U_\alpha P = P U_\alpha = U, \text{ and } \frac{1}{n} \sum_{k=0}^{n-1} U_\alpha^k \xrightarrow{so} P.$$

2. *If α is multiplicative, then*

(i) U_α is an isometry.

(ii) $U_\alpha U_\alpha^*$ is the orthogonal projection onto $\overline{\pi_\phi(\alpha(\mathfrak{A}))\xi_\phi}$, and thus, belongs to the commutant of $\pi_\phi(\alpha(\mathfrak{A}))$.

(iii) For all $a \in \mathfrak{A}$, $U_\alpha \circ \pi_\phi(a) = \pi_\phi(\alpha(a)) \circ T_\alpha$, and $\pi_\phi(a) = U_\alpha^* \pi_\phi(\alpha(a)) U_\alpha$, $U_\alpha \pi_\phi(a) U_\alpha^* = \pi_\phi(\alpha(a)) U_\alpha U_\alpha^*$.

A linear map between C^* -algebras, $\alpha : \mathfrak{A} \rightarrow \mathfrak{B}$ is called a *Schwarz map* if $\alpha(x)^*\alpha(x) \leq \alpha(x^*x)$. And if $\alpha : \mathfrak{A} \rightarrow \mathfrak{B}$ is positive, then for any $x \in \mathfrak{A}$,

$$\alpha(x)^*\alpha(x) + \alpha(x)\alpha(x)^* \leq \|\alpha\| \alpha(x^*x + xx^*).$$

Thus, the inequality in 1. of the lemma is satisfied by Schwarz maps, or when ϕ is tracial and $\|\alpha\| \leq 1$.

Proof. 1.

(i) U_α is a contraction follows from the same calculation as lemma 5.1.8, where we have an inequality instead of the third equality.

(ii) Let $\{u_\lambda\}$ be a approximate identity for \mathfrak{A} . Then,

$$\begin{aligned} \|\xi_\phi - U_\alpha(\widehat{u_\lambda})\| &= \|\xi_\phi\|^2 + \|U_\alpha(\widehat{u_\lambda})\|^2 - 2 \operatorname{Re} \langle \xi_\phi, \widehat{u_\lambda} \rangle \\ &= \|\phi\| + \phi(\alpha(u_\lambda)^*\alpha(u_\lambda)) - 2 \operatorname{Re} \phi(\alpha(u_\lambda)) \\ &\leq 1 + \phi(u_\lambda^* u_\lambda) - 2 \operatorname{Re} \phi(u_\lambda) \rightarrow 0. \end{aligned}$$

Thus, $U_\alpha \widehat{u_\lambda} \rightarrow \xi_\phi$, and a similar computation shows $\widehat{u_\lambda} \rightarrow \xi_\phi$. Hence, $U_\alpha \xi_\phi = \xi_\phi$.

(iii) Clearly, $U_\alpha P = P$. By lemma 2.3.1, $\{\xi \in H_\phi \mid U_\alpha \xi = \xi\} = \{\xi \in H_\phi \mid U_\alpha^* \xi = \xi\}$. We have $U_\alpha^* P = P$, which gives for all $\xi, \eta \in H_\phi$, $\langle P U_\alpha \xi, \eta \rangle = \langle \xi, U_\alpha^* P \eta \rangle = \langle P \xi, \eta \rangle \implies P U_\alpha = P$. The final statement is the mean ergodic theorem.

2.

(i) It is lemma 5.1.8.

(ii) We know for any $T \in \mathcal{B}(H)$, $\overline{\operatorname{Im}(T)} = \ker(T^*)^\perp$, and if T is an isometry, then its image is closed. This gives $\operatorname{Im}(U_\alpha U_\alpha^*) = \operatorname{Im}(U_\alpha) = U_\alpha(\overline{\pi_\phi(\mathfrak{A})\xi_\phi}) = \overline{\pi_\phi(\alpha(\mathfrak{A})\xi_\phi)}$. And since U_α is a isometry, ie, $U_\alpha^* U_\alpha = \operatorname{Id}$, which implies $U_\alpha U_\alpha^*$ is a projection. The second statement follows from lemma 4.2.23.

(iii) For $x \in \mathfrak{A}$,

$$U_\alpha \pi_\phi(a) \pi_\phi(x) \xi_\phi = \pi_\phi(\alpha(ax)) \xi_\phi = \pi_\phi(\alpha(a)) \pi_\phi(\alpha(x)) \xi_\phi = \pi_\phi(\alpha(a)) U_\alpha(x) \xi_\phi.$$

Taking adjoint in the equation, we get $\pi_\phi(a) U_\alpha^* = U_\alpha^* \pi_\phi(\alpha(a))$. The last results follow using $U_\alpha^* U_\alpha = 1_{\mathcal{H}_\phi}$.

□

Theorem 5.1.12. *Let \mathfrak{A} be a C^* -algebra, ϕ a state and $\alpha : \mathfrak{A} \rightarrow \mathfrak{A}$ a positive linear map such that $\phi \circ \alpha = \phi$ and for all $x \in \mathfrak{A}$,*

$$\phi(\alpha(x)^* \alpha(x)) \leq \phi(x^* x).$$

Then, for every $x \in \mathfrak{A}$ and every $\varepsilon > 0$, there is a relatively dense set $N \subset \mathbb{N}$ such that, for all $n \in n$,

$$\operatorname{Re} \phi(\alpha^n(x^*) x) \geq |\phi(x)|^2 - \varepsilon.$$

Proof. From the lemma, there exists $n_0 \in \mathbb{N}$ such that $\left\| \frac{1}{n_0} \sum_{k=0}^{n_0} U_\alpha^k \hat{x} - P\hat{x} \right\| \leq \frac{\varepsilon}{\|\xi\|}$. For $n \geq n_0$, since U_α is a contraction,

$$\begin{aligned} \left\| \frac{1}{(n_0+1)^2} \sum_{j,k=0}^{n_0} U^{l+k+l-j} \hat{x} - P\hat{x} \right\| &= \left\| \frac{1}{n_0+1} \sum_{j=0}^{n_0} U^{l-j} \left(\frac{1}{n_0+1} \sum_{k=0}^{n_0} U^k \hat{x} - P\hat{x} \right) \right\| \\ &\leq \left\| \left(\frac{1}{n_0+1} \sum_{k=0}^{n_0} U^k \hat{x} - P\hat{x} \right) \right\| \leq \frac{\varepsilon}{\|\xi\|}. \end{aligned}$$

Thus,

$$\left| \left\langle \frac{1}{(n_0+1)^2} \sum_{j,k=0}^{n_0} U^{l+k+l-j} \hat{x} - P\hat{x}, \hat{x} \right\rangle \right| \leq \varepsilon.$$

We also have,

$$\langle P\hat{x}, \xi_\phi \rangle = \lim_{n \rightarrow \infty} \left\langle \frac{1}{n} \sum_{k=0}^{n-1} U_\alpha^k \hat{x}, \xi_\phi \right\rangle = \lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} \phi(\alpha^k(x)) = \phi(x).$$

Hence, $|\phi(x)| \leq \|P\hat{x}\|$. Finally, note $\phi(x^* \alpha(x)) = \langle \pi_\phi(x^* \alpha(x)) \xi_\phi, \xi_\phi \rangle = \langle U_\alpha(\hat{x}), \hat{x} \rangle$, which gives

$$\begin{aligned} \frac{1}{(n_0+1)^2} \sum_{j,k=0}^{n_0} \operatorname{Re} \phi(x^* \alpha^{l+k-j}(x)) &= \operatorname{Re} \frac{1}{(n_0+1)^2} \sum_{j,k=0}^{n_0} \langle U_\alpha^{l+k-j} \hat{x}, \hat{x} \rangle \\ &= \operatorname{Re} \left(\frac{1}{(n_0+1)^2} \sum_{j,k=0}^{n_0} \langle U_\alpha^{l+k-j} \hat{x} - P\hat{x}, \hat{x} \rangle \right) + \operatorname{Re} \langle P\hat{x}, \hat{x} \rangle \geq |\phi(x)|^2 - \varepsilon. \end{aligned}$$

Therefore, for all $l \in \mathbb{N}$, there exists some $j_l, k_l \in [0, n_0]$ such that $\operatorname{Re} \phi(x^* \alpha^{l+k_l-j_l}(x)) \geq |\phi(x)|^2 - \varepsilon$. Since, $(l-1)n_0 \leq ln_0 + k_l - j_l \leq (l+1)n_0$, we have intervals of length $2n_0$, say $[(l-1)n_0, (l+1)n_0]$, will contain $n = ln_0 + k_l - j_l$ such that $\operatorname{Re} \phi(\alpha^n(x^*) x) \geq |\phi(x)|^2 - \varepsilon$. □

Corollary 5.1.13. *Let \mathfrak{A} be a C^* -algebra, ϕ a state and $\alpha : \mathfrak{A} \rightarrow \mathfrak{A}$ a positive linear map such that $\phi \circ \alpha = \phi$ and for all $x \in \mathfrak{A}$,*

$$\phi(\alpha(x)^* \alpha(x)) \leq \phi(x^* x).$$

Then, for every $a \in \mathfrak{A}$ with $\phi(a) > 0$,

$$\liminf_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} |\phi(\alpha^k(a^*) a)| > 0.$$

Proof. Same as theorem 5.1.2. \square

We extend another recurrence result.

Definition 5.1.14 (IP-set). $N \subset \mathbb{N}$ is called *IP-set* if there exists $p_1, p_2, \dots \in \mathbb{N}$ which generate N , that is,

$$N = \{p_{j_1} + p_{j_2} + p_{j_3} + \dots + p_{j_n} \mid 1 \leq j_1 < j_2 < \dots < j_n, n \in \mathbb{N}\}.$$

Theorem 5.1.15. Let \mathfrak{A} be a C^* -algebra, ϕ a state and $\alpha : \mathfrak{A} \rightarrow \mathfrak{A}$ a $*$ -homomorphism such that $\phi \circ \alpha = \phi$. Then, for all $x \in \mathfrak{A}$, IP-set $N \subset \mathbb{N}$ and $\varepsilon > 0$, the set

$$\{n \in N \mid \operatorname{Re} \phi(x^* \alpha^n(x)) \geq |\phi(x)|^2 - \varepsilon\} \text{ is infinite.}$$

Proof. Let N be generated by $p_1, p_2, \dots \in \mathbb{N}$. Denote $n_k = \sum_{j=1}^{j_k} p_j$ and choose a sequence $1 \leq j_1 < j_2 < \dots$ such that for all $k \in \mathbb{N}$, $p_{k+1} > n_k$. So, $i < j \implies n_j - n_i \in N$ and $i' < j', i < j < j' \implies n_{j'} - n_{i'} > n_j - n_i$.

Now, assume the set $\{n \in N \mid \operatorname{Re} \phi(x^* \alpha^n(x)) \geq |\phi(x)|^2 - \varepsilon\}$ is finite. Then, there is some $K \in \mathbb{N}$ such that

$$\operatorname{Re} \phi(\alpha^{n_l}(x)^* \alpha^{n_k}(x)) = \operatorname{Re} \phi(x^* \alpha^{n_k - n_l}(x)) < |\phi(x)|^2 - \varepsilon, \text{ for } k > l \geq K.$$

Note, for $n, m \in \mathbb{N}$,

$$\begin{aligned} & \langle \pi_\phi(\alpha^n(x))\xi_\phi - \phi(x)\xi_\phi, \pi_\phi(\alpha^m(x))\xi_\phi - \phi(x)\xi_\phi \rangle \\ &= \phi(\alpha^m(x)^* \alpha^n(x)) - \overline{\phi(x)}\phi(\alpha^n(x)) - \phi(x)\overline{\phi(\alpha^m(x))} + |\phi(x)|^2 \\ &= \phi(\alpha^m(x)^* \alpha^n(x)) - |\phi(x)|^2. \end{aligned}$$

Thus, for all $m > K$,

$$\begin{aligned} & \left\| \sum_{k=K}^{K+m-1} (\pi_\phi(\alpha^{n_k}(x))\xi_\phi - \phi(x)\xi_\phi) \right\|^2 \\ &= \sum_{k,l=K}^{K+m-1} \langle \pi_\phi(\alpha^{n_k}(x))\xi_\phi - \phi(x)\xi_\phi, \pi_\phi(\alpha^{n_l}(x))\xi_\phi - \phi(x)\xi_\phi \rangle \\ &= \sum_{k,l=K}^{K+m-1} (\phi(\alpha^{n_l}(x)^* \alpha^{n_k}(x)) - |\phi(x)|^2) \\ &= m\phi(x^*x) + \sum_{l,k=K, l < k}^{K+m-1} \operatorname{Re} \phi(\alpha^{n_l}(x)^* \alpha^{n_k}(x)) - m^2|\phi(x)|^2 \\ &\leq m\phi(x^*x) + m(m-1)(|\phi(x)|^2 - \varepsilon) - m^2|\phi(x)|^2 \rightarrow -\infty \text{ as } m \rightarrow \infty. \end{aligned}$$

This is a contradiction. In the last line, we use the fact that the number of pairs of integers (l, k) , $l \neq k$, $l, k < m$ is $m(m-1)$. \square

We end the section with the Poincare recurrence for von Neumann algebras.

Theorem 5.1.16. Let \mathcal{M} be a von Neumann algebra, ϕ a faithful normal state on \mathcal{M} , and $\alpha : \mathcal{M} \rightarrow \mathcal{M}$ be a $*$ -homomorphism such that $\phi \circ \alpha = \phi$. Then, for every $p \in \mathcal{M}$ and every $n \in \mathbb{N}$,

$$\bigvee_{k=n}^{\infty} \alpha^k(p) = \bigvee_{k=0}^{\infty} \alpha^k(p) \geq p.$$

This implies Poincare recurrence: for every projection $p \in \mathcal{M}$,

$$p \wedge \bigwedge_{n=0}^{\infty} \bigvee_{k=n}^{\infty} \alpha^k(p) = p.$$

Proof. We only need to prove $\bigvee_{k=n}^{\infty} \alpha^k(p) = \bigvee_{k=0}^{\infty} \alpha^k(p)$, ie, $\bigvee_{k=n}^{\infty} \alpha^k(p)$ is a constant, independent of n . We start by showing α is unital and normal.

Unital: Clearly, $\alpha(1_{\mathcal{M}})$ is a projection and hence, $\alpha(1_{\mathcal{M}}) \leq 1_{\mathcal{M}}$. Now, $\phi(1_{\mathcal{M}} - \alpha(1_{\mathcal{M}})) = 0$, which by faithfulness of ϕ gives $\alpha(1_{\mathcal{M}}) = 1_{\mathcal{M}}$.

Normal: Consider an increasing net $0 \leq a_i \uparrow a \in \mathcal{M}$. Then, $\alpha(a_i)$ is an increasing net, bounded by $\alpha(a)$. And, $\lim_i \alpha(a_i) = \alpha(a)$ for some $b \in \mathcal{M}$. Normality of ϕ gives,

$$\phi(b) = \lim \phi(\alpha(a_i)) = \lim \phi(a_i) = \phi(a) = \phi(\alpha(a)).$$

Since $\alpha(a) - b$ is positive, by faithfulness of ϕ , $\alpha(a) = b$.

Now, since $\ker(\alpha)$ is a w-closed two-sided ideal of \mathcal{M} , there is a central projection $q \in \mathcal{M}$ such that $\ker(\alpha) = \mathcal{M}(1_{\mathcal{M}} - q)$. Since α is w-continuous, $\alpha(\mathcal{M})$ is a von Neumann algebra, and thus, $\mathcal{M}q \cong \alpha(\mathcal{M})$. Finally,

$$\begin{aligned} \alpha \left(\bigvee_{k=n}^{\infty} \alpha^k(p) \right) &= \alpha \left(\left(\bigvee_{k=n}^{\infty} \alpha^k(p) \right) q \right) \quad (\text{for any } x \in \mathcal{M}, \alpha(x - xq) = 0) \\ &= \alpha \left(\left(\bigvee_{k=n}^{\infty} \alpha^k(p)q \right) \right) \quad (\text{follows from } q \text{ being central}) \\ &= \bigvee_{k=n}^{\infty} \alpha(\alpha^k(p)q) \\ &= \bigvee_{k=n}^{\infty} \alpha^{k+1}(p) \end{aligned}$$

The penultimate equality follows from normality of α and the fact that for any projections e_1, e_2 in a von Neumann algebras and a *-isomorphism $\mathcal{N}, \Phi : \mathcal{N} \rightarrow \mathcal{N}, \Phi(e_1 \vee e_2) = \Phi(e_1) \vee \Phi(e_2)$. This can be proved using $e_1 \vee e_2$ being the smallest projection containing e_1 and e_2 .

Applying ϕ to the previous result and using $\phi \circ \alpha = \alpha$, we get $\phi(\bigvee_{k=n}^{\infty} \alpha^k(p))$ is constant. Faithfulness of ϕ gives $\bigvee_{k=n}^{\infty} \alpha^k(p)$ is also constant. \square

5.2 Multiple recurrence in compact systems

Definition 5.2.1 (Almost periodic systems). A C*-dynamical system $(\mathfrak{A}, \phi, \alpha)$ is *almost periodic* if $\{U^n(\hat{x}) \mid n \in \mathbb{N}\}$ is precompact in H_ϕ for all $x \in \mathfrak{A}$.

Theorem 5.2.2. Let $(\mathfrak{A}, \phi, \alpha)$ be an almost periodic state preserving C*-dynamical system, with the support projection of ϕ in the double dual \mathfrak{A}^{**} being central, then for every $p \in \mathbb{N}, m_1, m_2, \dots, m_p \in \mathbb{N}, x_1, x_2, \dots, x_p \in \mathfrak{A}$ and $\varepsilon > 0$, there is a relatively dense set $N \subset \mathbb{N}$ such that

$$|\phi(\alpha^{m_1 n}(x_1)\alpha^{m_2 n}(x_2)\dots\alpha^{m_p n}(x_p)) - \phi(x_1 x_2 \dots x_p)| \leq \varepsilon \text{ for all } n \in N,$$

It follows that

$$\phi(\alpha^{m_1 n}(x_1)\alpha^{m_2 n}(x_2)\dots\alpha^{m_p n}(x_p)) \geq \phi(x_1 x_2 \dots x_p) - \varepsilon \text{ for all } \varepsilon > 0.$$

As in the proof of theorem 5.1.5., we have the multiple recurrence property: for $0 < a \in \mathfrak{A}$, if $\phi(a) > 0$, then

$$\liminf_{n \rightarrow \infty} \frac{1}{n+1} \sum_{k=0}^n |\phi(a \alpha^{m_1 k}(a) \dots \alpha^{m_p n}(a))| > 0.$$

For the proof, we start with some results from metric spaces.

Lemma 5.2.3. *Let (X, d) is a totally bounded metric space and $\epsilon > 0$, then the set $\{n \in \mathbb{N} \mid \text{there are } x_1, \dots, x_n \in X \text{ such that } d(x_j, x_k) > \epsilon \text{ for } j \neq k\}$ is bounded.*

Proof. Assuming the contrary, we have for any $n \in \mathbb{N}$, $\omega_1^{(1)}, \dots, \omega_n^{(n)} \in \Omega$ such that $d(\omega_j^{(n)}, \omega_k^{(n)}) > \epsilon$, $j \neq k$. Consider the completion of X , \tilde{X} , which will be a compact metric space. Now, let $\lim_{j \leq n \rightarrow \infty} \omega_j^n = \tilde{\omega}_j$. Then, $d(\tilde{\omega}_j, \tilde{\omega}_k) = \lim_{k \leq n \rightarrow \infty} d(\omega_j^{(n)}, \omega_k^{(n)}) \geq \epsilon$ when $j < k$. So, the sequence $\tilde{\omega}_j$ cannot have any Cauchy subsequence in $\tilde{\Omega}$, which contradicts its compactness. \square

Corollary 5.2.4. *Let (Ω, d) be a metric space, $T : \Omega \rightarrow \Omega$ an isometry, and $\omega \in \Omega$, then the following statements are equivalent:*

- (i) *the orbit of ω , $\{T^n \omega \mid n \in \mathbb{N}\}$ is totally bounded.*
- (ii) *for all $\epsilon > 0$ there exists a relatively dense set $N \subset \mathbb{N}$ such that $d(T^n \omega, \omega) \leq \epsilon$, for all $n \in N$.*

Proof. (ii) \implies (i): For $\epsilon > 0$, let $L > 0$ be the length of the interval such that any interval of length L contains a $n \in N$ such that $d(T^n \omega, \omega) \leq \epsilon$. We show $\{B_\epsilon(T^j \omega) \mid 1 \leq j \leq L\}$ cover the orbit of ω : for any $k \in \mathbb{N}$, there is a $1 \leq j \leq L$ such that $k - j \in N$. Then $d(T^k \omega, T^j \omega) \leq d(T^{k-j} \omega, \omega) \leq \epsilon$.

- (i) \implies (ii): From the previous lemma, let p be the largest integer such that there are $n_1 < \dots < n_p \in \mathbb{N}$ such that $d(T^{n_j} \omega, T^{n_k} \omega) > \epsilon$, $j \neq k$. Now, consider any interval of length n_p , that is, $[a, a + n_p]$ for some $a \in \mathbb{N}$. Then $\{a + n_j\}_{j=1}^p$ are also integers such that

$$d(T^{a+n_j} \omega, T^{a+n_k} \omega) = d(T^{n_j} \omega, T^{n_k} \omega) > \epsilon, \quad j \neq k.$$

Hence, by maximality of p , it must be that $d(T^{a+n_j} \omega, \omega) \leq \epsilon$ for some $1 \leq j \leq p$. \square

For $1 \leq i \leq p$, if (Ω_i, d_i) are metric spaces, $T_i : \Omega_i \rightarrow \Omega_i$ isometries and $\omega_i \in \Omega_i$, then all orbits $T_i^n(\omega_i)$, $1 \leq i \leq p$ are totally bounded if and only if for every $\epsilon > 0$, there exists a relatively dense set $N \subset \mathbb{N}$ such that $d_i(T_i^n \omega_i, \omega_i)$, $1 \leq i \leq p$, $n \in N$. This follows from the above corollary with $\Omega = \prod \Omega_i$, $d((\omega_i), (\rho_i)) = \max_i d_i(\omega_i, \rho_i)$, $T = \prod T_i$.

We now use this in the following way. Consider a contraction U on a Hilbert space H . Define the set of almost periodic vectors,

$$H_{AP}^U = \{\xi \in H \mid \{U^n(\xi) \mid n \in \mathbb{N}\} \text{ is relatively norm-compact}\}.$$

It is clearly a U -invariant, linear subspace of H and it is easy to show that it closed.

Theorem 5.2.5. *Let \mathfrak{A} be a C^* -algebra, ϕ a state on \mathfrak{A} such that support $s(\phi)$ is central in \mathfrak{A}^{**} is central, and $\alpha : \mathfrak{A} \rightarrow \mathfrak{A}$ a positive linear map such that $\phi \circ \alpha = \phi$.*

- (i) *If $\phi(\alpha(x)^* \alpha(x)) \leq \phi(x^* x)$, then there is a normal positive linear map, $\Psi : \pi_\phi(\mathfrak{A})'' \rightarrow \pi_\phi(\mathfrak{A})''$ such that $\Psi(T)\xi_\phi = U_\alpha(T\xi_\phi)$ for $T \in \pi_\phi(\mathfrak{A})''$, $\Psi(1_{H_\phi}) = 1_{H_\phi}$, $\|\Psi\| \leq \|\alpha\|$, preserving $\omega_{\xi_\phi} \mid \pi_\phi(\mathfrak{A})''$ and $\Psi(\pi_\phi(a)) = \pi_\phi(\alpha^{**}(a))$, $a \in \mathfrak{A}^{**}$.*
- (ii) *If α is multiplicative, then $\pi_\phi(\alpha(\mathfrak{A}))''$ is a von Neumann subalgebra of $\pi_\phi(\mathfrak{A})''$ and the central support of the projection $U_\alpha U_\alpha^*$ in $\pi_\phi(\alpha(\mathfrak{A}))''$ is 1_{H_ϕ} . And, $\Psi(T)U_\alpha U_\alpha^* = U_\alpha T U_\alpha^*$.*

Proof. (i) Let $T \in \pi_\phi(\mathfrak{A})''$. By Kaplansky density theorem, there is a net $(a_\iota) \in \mathfrak{A}$ such that $\|a_\iota\| \leq \|T\|$ and $\pi_\phi(a_\iota) \xrightarrow{\text{so}} T$. Define

$$\pi_\phi(\alpha(a_\iota)) \xrightarrow{\text{wo}} \Psi(T) \in \pi_\phi(\mathfrak{A})''.$$

Then, $\pi_\phi(\alpha(a_\iota))\xi_\phi \xrightarrow{\text{weak}} \Psi(T)\xi_\phi$ using Riesz representation theorem, and also $\pi_\phi(\alpha(a_\iota))\xi_\phi = U_\alpha\pi_\phi(a_\iota)\xi_\phi \rightarrow U_\alpha(T\xi_\phi)$. Hence, $\Psi(T)\xi_\phi = U_\alpha(T\xi_\phi)$. Using the fact that ξ_ϕ is a cyclic vector for $\pi_\phi(\mathfrak{A})'$, it is easy to see that $\Psi(T)$ is uniquely determined by the value $\Psi(T)\xi_\phi = U(T\xi_\phi)$ and Ψ is a well-defined map. Hence, $\Psi(T) = U_\alpha T$ and linearity of the map Ψ is clear.

For w-continuity, by $\pi_\phi(\mathfrak{A})'\xi_\phi = H_\phi$, it is enough to show the linear functionals, $\omega_{T'_1\xi_\phi, T'_2\xi_\phi} \circ \Psi$ for all $T'_1, T'_2 \in \mathfrak{A}'$ are w-continuous:

$$\omega_{T'_1\xi_\phi, T'_2\xi_\phi} \circ \Psi(T) = \langle \Psi(T)T'_1\xi_\phi, T'_2\xi_\phi \rangle = \langle \Psi(T)\xi_\phi, (T'_1)^*T'_2\xi_\phi \rangle = \langle T, U_\alpha^*(T'_1)^*T'_2\xi_\phi \rangle.$$

And $\|\Psi\| \leq \|\alpha\|$:

$$\|\Psi(T)\| = \sup_{\xi, \eta \in (H_\phi)_1} \langle \Psi(T)\xi, \eta \rangle = \sup_{\xi, \eta \in (H_\phi)_1} \lim_\iota \langle \pi_\phi(\alpha(a_\iota))\xi, \eta \rangle \leq \lim_\iota \|\alpha\| \|a_\iota\| \leq \|\alpha\| \|T\|.$$

We have $\Psi(1_{H_\phi})\xi_\phi = U_\alpha(\xi_\phi) = \xi_\phi$, which gives $\Psi(1_{H_\phi}) = 1_{H_\phi}$. Also,

$$\omega_{\xi_\phi}(\Psi(T)) = \langle \Psi(T)\xi_\phi, \xi_\phi \rangle = \langle U_\alpha(T\xi_\phi), \xi_\phi \rangle = \langle T\xi_\phi, U_\alpha^*\xi_\phi \rangle = \omega_{\xi_\phi}(T).$$

The positivity of Ψ follows with the observation $\Psi(\pi_\phi(a)) = \pi_\phi(\alpha(a))$, $\alpha \in \mathfrak{A}$ and the final assertion follows from theorem 4.2.41 and density of \mathfrak{A} in its double dual.

(ii) Since $\pi_\phi(\mathfrak{A})$ is w-dense in $\pi_\phi(\mathfrak{A})''$, we have $\pi_\phi(a_\iota) \xrightarrow{w} 1_{H_\phi}$, and since $\Psi(1_{H_\phi}) = 1_{H_\phi}$, by w-continuity of Ψ , we have that 1_{H_ϕ} is in the w-closure $\pi_\phi(\mathfrak{A})$. By bicommutant theorem, $\pi_\phi(\mathfrak{A})$ is a von Neumann algebra.

Call the central support of $U_\alpha U_\alpha^*$, $P \in \pi_\phi(\alpha(\mathfrak{A}))''$. Recall, $U_\alpha U_\alpha^*$ is the orthogonal projection onto $\pi_\phi(\alpha(\mathfrak{A}))\xi_\phi$, to which 1_{H_ϕ} belongs. Thus, P fixes ξ_ϕ and since ξ_ϕ is separating for $\pi_\phi(\mathfrak{A})''$, we get $P = 1_{H_\phi}$. \square

Lemma 5.2.6. *Let \mathcal{H} be a Hilbert space, $\mathcal{M} \subset \mathcal{B}(\mathcal{H})$ a closed linear subspace, $\xi_0 \in H$ such that $\overline{\mathcal{M}'\xi_0} = H$ and U a linear isometry on H such that $\mathcal{M}\xi_0 \subset H_{AP}^U$. Let*

$$\mathcal{G} = \{\text{linear contractions } \theta : H \rightarrow H \mid \theta(T)\xi_0 \in \overline{\{U^n T\xi_0 \mid n \in \mathbb{N}\}} \text{ for all } T \in \mathcal{M}\}.$$

For any $p \in \mathbb{N}$, let $\theta_1, \dots, \theta_p \in \mathcal{G}$, then,

For any integer $p \in \mathbb{N}$, $\theta_1, \dots, \theta_p \in \mathcal{G}$, $T_1, \dots, T_p \in \mathcal{A}$, $\xi_1, \dots, \xi_p \in \mathcal{H}$ and $\varepsilon > 0$, there exists a relatively dense set $N \subset \mathbb{N}$ such that

$$\|\theta_j^n(T_j)\xi_j - T_j\xi_j\| \leq \varepsilon \text{ for all } 1 \leq j \leq p, n \in N.$$

Proof. Note,

$$\text{if } \xi \in \overline{\{U^n(T\xi_0) \mid n \in \mathbb{N}\}}, \text{ then } \|\xi\| = \|\xi_0\|.$$

Hence, for $\theta \in \mathcal{G}$, by assumption, we define $\|\theta(T)\xi\| = \|\xi_0\|$ and we define a linear isometry

$$U_\theta : \overline{\mathcal{M}\xi_0} \rightarrow \overline{\mathcal{M}\xi_0}, U_\theta((T\xi_0)) = \theta(T)\xi_0.$$

Then, for $\theta_1, \theta_2 \in \mathcal{G}$, $U_{\theta_1 \circ \theta_2} = U_{\theta_1} U_{\theta_2}$ and for $\theta \in \mathcal{G}$, $n \in \mathbb{N}$, $U_{\theta^n} = U_\theta^n$. Thus, for $T \in \mathcal{M}$, $\{U_\theta^n(T\xi_0) \mid n \in \mathbb{N}\} = \{\theta^n(T)\xi_0 \mid n \in \mathbb{N}\} \subset \{U^n(T)\xi_0 \mid n \in \mathbb{N}\}$, which is relatively compact, and this implies that $\mathcal{M}\xi_0 \subset (\overline{\mathcal{M}\xi_0})_{AP}^{U_\theta}$.

Choose $T'_1, \dots, T'_p \in \mathcal{M}'$ such that $\|\xi_j - T'_j \xi_0\| \leq \varepsilon/3 \|T_j\|$, $1 \leq j \leq p$. By Corollary 5.2.4., with $T_j = U_{\theta_j}$, we get a relatively dense set $N \subset \mathbb{N}$ such that

$$\|\theta_j^n(T_j)\xi_0 - T_j\xi_0\| = \|U_{\theta_j}^n(T_j\xi_0) - T_j\xi_0\| \leq \frac{\varepsilon}{3\|T'_j\|}, \quad 1 \leq j \leq p, \quad n \in N.$$

Finally, for $1 \leq j \leq p, n \in N$,

$$\begin{aligned} \|\theta_j^n(T_j)\xi_j - T_j\xi_j\| &\leq \|\theta_j^n(T_j)\xi_j - \theta_j^n(T_j)(T'_j\xi_0) + \theta_j^n(T_j)(T'_j\xi_0) - T'_jT_j\xi_0 + T'_jT_j\xi_0 - T_j\xi_j\| \\ &\leq \|\theta_j^n(T_j)(\xi_j - T'_j\xi_0)\| + \|T'_j(\theta_j^n(T_j)\xi_0 - T_j\xi_0)\| + \|T_j(T'_j\xi_0 - \xi_j)\| \leq \varepsilon. \end{aligned}$$

□

Theorem 5.2.7. *Let \mathfrak{A} be a C^* -algebra, ϕ a state on \mathfrak{A} such that support $s(\phi)$ in \mathfrak{A}^{**} is central, and let $\mathcal{M}_{AP} = \{T \in \pi_\phi(\mathfrak{A})'' \mid T\xi_\phi \in (H_\phi)_{AP}\}$. Then, for any $p \in \mathbb{N}$, $m_1, \dots, m_p \geq 1, T_1, \dots, T_p \in \mathcal{M}_{AP}$, $\xi_1, \dots, \xi_p \in H_\phi$ and $\varepsilon > 0$, there exists a relatively dense set $N \subset \mathbb{N}$ such that*

$$\|\Psi^{m_1 n}(T_1)\xi_1 - T_1\xi_1\| \leq \varepsilon \text{ for all } 1 \leq j \leq p, \quad n \in N.$$

Proof. Note, $\Psi(T)\xi_\phi = U(T\xi_\phi)$ implies $\Psi|\mathcal{M}_{AP} \in \mathcal{G}$. We apply the above lemma with $\mathcal{M} = \mathcal{M}_{AP}$, $\theta_j = (\Psi|\mathcal{M}_{AP})^{m_j}$ □

Corollary 5.2.8. *Let $(\mathfrak{A}, \phi, \alpha)$ be a state-preserving C^* -dynamical system such that $s(\phi)$ in \mathfrak{A}^{**} is central. Then, for any $p \in \mathbb{N}$, $m_1, \dots, m_p \geq 1, T_1, \dots, T_p \in \mathcal{M}_{AP}$, $\xi_1, \dots, \xi_p \in H_\phi$ and $\varepsilon > 0$, there exists a relatively dense set $N \subset \mathbb{N}$ such that*

$$\|\Psi^{m_1 n}(T_1)\Psi^{m_2 n}(T_2) \dots \Psi^{m_p n}(T_p)\xi - T_1T_2 \dots T_p\xi\| \leq \varepsilon.$$

Proof. Apply above theorem with $\xi_j = T_{j+1} \dots T_p\xi$, we get a relatively dense set $N \subset \mathbb{N}$ such that for all $1 \leq j \leq p$ and $n \in N$,

$$\|(\Psi^{m_j n}(T_j) - T_j)T_{j+1} \dots T_p\xi\| \leq \varepsilon(p\|T_1\| \|T_{j-1}\|)^{-1}.$$

Using the identity,

$$\prod_{i=1}^k a_i - \prod_{i=1}^k b_i = \sum_{j=1}^k \left[\left(\prod_{i=1}^j a_i \right) (a_j - b_j) \left(\prod_{i=j+1}^k b_i \right) \right],$$

we get,

$$\begin{aligned} &\|\Psi^{m_1 n}(T_1)\Psi^{m_2 n}(T_2) \dots \Psi^{m_p n}(T_p)\xi - T_1T_2 \dots T_p\xi\| \\ &\leq \left\| \sum_{j=1}^p \Psi^{m_1 n}(T_1) \dots \Psi^{m_{j-1} n}(T_{j-1}) \dots \Psi^{m_{j-1} n}(T_{j-1})(\Psi^{m_j n}(T_j) - T_j)S_j T_{j+1} \dots T_p \right\| \\ &\leq \sum_{j=1}^p \|T_1\| \dots \|T_{j-1}\| \|(\Psi^{m_j n}(T_j) - T_j)T_{j+1} \dots T_p\xi\| \leq \varepsilon. \end{aligned}$$

□

We now prove theorem 5.2.2.

Proof of theorem 5.2.2. Since the system is almost periodic, we have $\pi_\phi(\mathfrak{A}) \subset \mathcal{M}_{AP}$ and we can apply previous corollary with $T = \pi_\phi(x)$ and $\xi = \xi_\phi$. Note $\Psi(\pi_\phi(x))\xi_\phi = U_\alpha(\pi_\phi(x)\xi_\phi) = \pi_\phi(\alpha(x))\xi_\phi$. We have a relatively dense set $N \subset \mathbb{N}$ such that for all $1 \leq j \leq p$ and $n \in N$,

$$\begin{aligned} & \|\pi_\phi(\alpha^{m_1 n}(x_1)\alpha^{m_2 n}(x_2)\dots\alpha^{m_p n}(x_p))\xi_\phi - \pi_\phi(x_1 x_2 \dots x_p)\xi_\phi\| \\ &= \|\Psi^{m_1 n}(T_1)\Psi^{m_2 n}(T_2)\dots\Psi^{m_p n}(T_p)\xi - T_1 T_2 \dots T_p \xi\| \leq \varepsilon. \end{aligned}$$

Finally,

$$\begin{aligned} & |\phi(\alpha^{m_1 n}(x_1)\alpha^{m_2 n}(x_2)\dots\alpha^{m_p n}(x_p)) - \phi(x_1 x_2 \dots x_p)| \\ &= |\langle \pi_\phi(\alpha^{m_1 n}(x_1)\alpha^{m_2 n}(x_2)\dots\alpha^{m_p n}(x_p))\xi_\phi - \pi_\phi(x_1 x_2 \dots x_p)\xi_\phi, \xi_\phi \rangle| \\ &\leq \|\pi_\phi(\alpha^{m_1 n}(x_1)\alpha^{m_2 n}(x_2)\dots\alpha^{m_p n}(x_p))\xi_\phi - \pi_\phi(x_1 x_2 \dots x_p)\xi_\phi\| \leq \varepsilon. \end{aligned}$$

□

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