

I.1 MOTIVATION1.1 LAW OF THE LARGE NUMBERS

$X_1, X_2, \dots : \Omega \rightarrow \mathbb{R}$ iid random variables ($\mathbb{E}[X_1] < \infty$),

then

• $\lim_{n \rightarrow \infty} \mathbb{P}\left[\left|\frac{1}{n} \sum_{k=1}^n X_k - \mathbb{E}[X_1]\right| > \varepsilon\right] = 0 \quad \forall \varepsilon > 0$ (weak)

• $\mathbb{P}\left[\lim_{n \rightarrow \infty} \frac{1}{n} \sum_{k=1}^n X_k = \mathbb{E}[X_1]\right] = 1$ (strong)

Is independence necessary? \rightarrow No!

1.2 A SIMPLE WEATHER MODEL

X_i : i^{th} day's weather $\in \{R, S\}$ "raining" / "sunny"

$\mathbb{P}[X_{i+1} = b \mid X_i = a] = P_{ab} \in \mathbb{R}^{2 \times 2}, \quad a, b \in \{R, S\}$

Let

$$P = \begin{bmatrix} 0.6 & 0.4 \\ 0.2 & 0.8 \end{bmatrix} \begin{matrix} R \\ S \end{matrix} \begin{matrix} R \\ S \end{matrix}$$

$\{X_i\}_{i \geq 0}$ is a Markov chain

$\mathbb{P}_S[X_i = S] = [0 \ 1] P^i \begin{bmatrix} 0 \\ 1 \end{bmatrix}$

\uparrow " $X_0 = S$ "

Eigendecomposition:

$M = \begin{bmatrix} 1/3 & 2/3 \\ -1/2 & 1/2 \end{bmatrix}$

$MPM^{-1} = \begin{bmatrix} 1 & \\ & 0.4 \end{bmatrix}$

v, w : left eigenvectors

v^T

w

λ_1

λ_2

λ_2

1.3 LONG-TERM BEHAVIOR

$$\mathbb{E} \left[\frac{\text{\# sunny days in the first } m \text{ days}}{m} \right] = \begin{matrix} =: W_m^S \text{ for sunny} \\ W_m^R \text{ for rainy} \end{matrix}$$

$$= \frac{1}{m} \sum_{i=0}^{m-1} \mathbb{P}_f [X_i = S] \quad f \in \mathbb{R}^2: \text{initial distribution}$$

$$= \frac{1}{m} \sum_{i=0}^{m-1} f^T P^i \begin{bmatrix} 0 \\ 1 \end{bmatrix}$$

$$= f^T M \frac{1}{m} \sum_{i=0}^{m-1} \begin{bmatrix} 1 & 0.4 \\ 0 & 0.4 \end{bmatrix}^i M \begin{bmatrix} 0 \\ 1 \end{bmatrix}$$

$$\xrightarrow{m \rightarrow \infty} f^T \begin{bmatrix} 1/3 & 2/3 \\ 1/3 & 2/3 \end{bmatrix} \begin{bmatrix} 0 \\ 1 \end{bmatrix} = 2/3 = N_2 = N_S = N_{\text{"sunny"}}$$

Observations:

(a) $\mathbb{E} [W_m^{S/R}] \xrightarrow{m \rightarrow \infty} N_{S/R}$

regardless the initial condition

(b) It holds even ("strong law"):

$$W_m^{S/R} \xrightarrow[m \rightarrow \infty]{\text{a.s.}} N_{S/R}$$

(c) $f \in \mathbb{R}_{\geq 0}^2, \|f\|_1 = 1: f^T P^m \xrightarrow{m \rightarrow \infty} N^T$ ("limiting distribution")

Is there something similar for deterministic systems as well?

1.4 DYNAMICAL SYSTEMS

X some set

Let $(G, +)$ be an (additive) group or semigroup, and

$(T_g)_{g \in G}$ a family of maps $T_g: X \rightarrow X$, such that

$$(i) T^{g_1} \circ T^{g_2} = T^{g_1+g_2} \quad \forall g_{1,2} \in G \quad | 3$$

$$(ii) T^0 = id$$

Then (X, T^\bullet) is called a dynamical system with state space X , and T^\bullet is called the flow.

- $G = \mathbb{N}_0$ or \mathbb{Z} : discrete time system, here $T := T^1$
- $G = \mathbb{R}_{\geq 0}$ or \mathbb{R} : continuous time system

1.5 THE STUDY OF DYNAMICAL SYSTEMS

(a) Differentiable dynamics: T differentiable, X smooth manifold
Local stretching ($\|DT\|$), Lyapunov exponents, ...
 \leadsto see e.g. [Ma]

(b) Topological dynamics: T continuous, X topological space
E.g. "transitivity": U, V open sets; is there a $n \geq 0$ such that
 $T^{-n}(U) \cap V \neq \emptyset$?
 \leadsto see e.g. [BS]

(c) Ergodic theory: T measurable, X measure space
E.g. long-term behavior as above

1.6 PREVIEW

Poincaré's Recurrence Theorem (PRT)

For almost all $x \in X$ (measured by the Lebesgue measure),
the system starting at x will return arbitrary close to x at arbitrary
large times.

Question: What other dynamical information can be extracted from measure-theoretic properties of a system?

The Ergodic Hypothesis

For certain measures μ , many functions $f: X \rightarrow \mathbb{R}$, and many states $x \in X$:

$$\lim_{t \rightarrow \infty} \frac{1}{t} \int_0^t f(T^s x) ds = \frac{1}{\mu(X)} \int_X f(x) d\mu(x) = "E_{\mu}[f]"$$

"time-average = space-average"; cf. law of the large numbers

Note: the ergodic hypothesis is a quantitative version of the PRT:

with $f = \chi_{B_{\epsilon}(x)}$ (characteristic function of an ϵ -ball)

we get the relative time the system spends in $B_{\epsilon}(x)$.

I.2 THE SETUP OF ERGODIC THEORY

1.7 MEASURE PRESERVING TRANSFORMATIONS

Def 1: A measure space is a triplet (X, \mathcal{B}, μ) where

(i) X is a set (also called "space")

(ii) \mathcal{B} is a σ -algebra (collection of subsets of X containing \emptyset and being closed under complements and countable unions)

$A \in \mathcal{B}$ is called measurable

(iii) $\mu: \mathcal{B} \rightarrow [0, \infty]$ is a measure, i.e. a σ -additive function on \mathcal{B} .

$\mu(X) = 1 \rightarrow$ "probability space"

Def B: A measure preserving transformation (mpt) is a quartet (X, \mathcal{B}, μ, T) where (X, \mathcal{B}, μ) is a measure space, and

(i) T is measurable, i.e. $E \in \mathcal{B} \Rightarrow T^{-1}E := T^{-1}(E) \in \mathcal{B}$

(ii) μ is invariant, i.e. $\mu(T^{-1}E) = \mu(E) \forall E \in \mathcal{B}$

A probability preserving transformation (ppt) is a mpt on a prob. space.

Ex:

$X = S^1$ (circle with unit circumference), $Tx = x + d \pmod{1}$, where $d \in \mathbb{R}$, μ is the Lebesgue measure

Def C: Two mpt's $(X_i, \mathcal{B}_i, \mu_i, T_i)$ are called **measure-theoretically isomorphic** if there is a measurable map $\pi: X_1 \rightarrow X_2$:

(i) there are $X_i' \in \mathcal{B}_i$ s.t. $\mu_i(X_i \setminus X_i') = 0$ and $\pi: X_1' \rightarrow X_2'$ is invertible with measurable inverse ("essentially a bijection")

(ii) for every $E \in \mathcal{B}_2$, $\pi^{-1}(E) \in \mathcal{B}_1$ and $\mu_1(\pi^{-1}E) = \mu_2(E)$

(iii) $T_2 \circ \pi = \pi \circ T_1$ on X_1

$$\begin{array}{ccc} X_1 & \xrightarrow{T_1} & X_1 \\ \pi \downarrow & & \downarrow \pi \\ X_2 & \xrightarrow{T_2} & X_2 \end{array} \quad \text{diagram commutes}$$

1.8 POINCARÉ'S RECURRENCE THEOREM

Thm: Let (X, \mathcal{B}, μ, T) be a ppt. If $E \in \mathcal{B}$ with $\mu(E) > 0$ then for almost every $x \in E$ there is a sequence $n_k \rightarrow \infty$ s.t. $T^{n_k}(x) \in E$.

Proof ([BG, Thm 3.2.1], extended)

Let $N_0 \subseteq E$ the set of points that never return to E , i.e.

$$N_0 := \{ x \in E \mid T^k(x) \notin E, \forall k \geq 1 \}.$$

(a) We show that $T^{-i}(N_0) \cap T^{-j}(N_0) = \emptyset \quad \forall i > j \geq 0$. If $x \in X$ s.t. there are $i > j$ with $T^i(x) \in N_0 \ni T^j(x)$, then

$$\underbrace{T^{i-j}(T^j(x))}_{=: y \in N_0} = T^i(x) \in N_0 \implies y \in N_0 \text{ is recurrent} \implies y \notin N_0 \quad \downarrow$$

(b) N_0 is measurable, since $E \in \mathcal{A}$, T measurable, and

$$N_0 = \bigcap_{k=1}^{\infty} T^{-k}(E^c)$$

From ppt assumption & (b) $\implies \mu(T^{-i}(N_0)) = \mu(N_0) \quad \forall i \geq 0$

Thus, from (a) \implies

$$1 = \mu(X) \geq \mu\left(\bigcup_{k=0}^{\infty} T^{-k}(N_0)\right) \stackrel{(a)}{=} \sum_{k=0}^{\infty} \mu(T^{-k}(N_0)) = \sum_{k=0}^{\infty} \mu(N_0)$$

Hence, $\mu(N_0) = 0$, and thus almost every point in E returns at least once to E : let $E_1 := E \setminus N_0$. Let N_1 be the set that never returns to E_1 , $i \geq 1$, and $E_{i+1} := E_i \setminus N_i$. We show iteratively that $\mu(N_i) = 0 \quad \forall i \geq 0$.

The recurrent set, we want, is

$$E_{\infty} := E \setminus \bigcup_{i=0}^{\infty} N_i,$$

clearly measurable, and

$$\mu(E) \geq \mu(E_{\infty}) = \mu(E) - \underbrace{\mu\left(\bigcup_{i=0}^{\infty} N_i\right)}_{< 0} \geq \mu(E) - \sum_{i=0}^{\infty} \mu(N_i)$$

1.9 REMARKS TO THE PRT

(a) The finiteness of μ is essential, as the proof shows.

E.g. take $X = \mathbb{Z}$, $\mathcal{B} = 2^X$, $T(x) = x+1$, and

$$\mu(E) = |E|, \text{ the counting measure}$$

Then, μ is invariant, but recurrence is clearly violated.

(b) $X = [0, 1]$, $T(x) = 10x \pmod{1}$

T leaves the Lebesgue measure invariant

PRT \Rightarrow for almost all $X \ni x = 0.\xi_1\xi_2\xi_3\dots\xi_m\dots$

the sequence $\xi_1\xi_2\dots\xi_m$ appears infinitely many times in the decimal expansion.

1.10 INVARIANT SETS

Def: $E \in \mathcal{B}$ is called **invariant**, if $T^{-1}(E) = E$.

Note: If E is invariant:

(i) so is $E^c := X \setminus E$;

(ii) $T: E \rightarrow E$ and $T: E^c \rightarrow E^c$ can be considered separately, they do not interact.

1.11 ERGODICITY

Def: A mpt (X, \mathcal{B}, μ, T) is **ergodic**, if every invariant set E satisfies $\mu(E) = 0$ or $\mu(E^c) = 0$. We call μ an ergodic measure.

"The system is essentially indecomposable"

Thm: For a mpt the following are equivalent:

- (a) μ is ergodic
- (b) if $E \in \mathcal{B}$ and $\mu(T^{-1}E \Delta E) = 0$, then $\mu(E) = 0$ or $\mu(E^c) = 0$
- (c) if $f: X \rightarrow \mathbb{R}$ is measurable and $f \circ T = f$ a.e., then $\exists c \in \mathbb{R}$ s.t. $f = c$ a.e.

Remarks:

- $A \Delta B := (A \setminus B) \cup (B \setminus A)$ is called the *symmetric difference*
- A set E with $\mu(T^{-1}E \Delta E) = 0$ is called *essentially T -invariant*
- Part (c) of the Thm can be read as "a.e. trajectory of T goes almost everywhere in X "

Proof:

"(a) \Rightarrow (b)": Let $E \in \mathcal{B}$ with $\mu(E \Delta T^{-1}E) = 0$. We'll construct an invariant $E_0 \in \mathcal{B}$ "close to E ". Let

$$E_0 = \{x \in X \mid T^k(x) \in E \text{ for infinitely many } k \in \mathbb{N}\}$$

E_0 is measurable and invariant (exercise). Now, to closeness:

- $x \in E_0 \setminus E$, then $\exists k: T^k(x) \in E$, but $x \notin E$, i.e. $x \in T^{-k}E \setminus E$
- $x \in E \setminus E_0$, then $\exists k: T^k(x) \notin E$, but $x \in E$, i.e. $x \in E \setminus T^{-k}E$

$$\Rightarrow E \Delta E_0 \subset \bigcup_{k \geq 1} E \Delta T^{-k}E. \text{ By } \mu(A \Delta C) \leq \mu(A \Delta B) + \mu(B \Delta C)$$

$$\mu(E \Delta E_0) \leq \mu\left(\bigcup_{k \geq 1} E \Delta T^{-k}E\right) \quad \text{⊗ (prove this!)}$$

$$\leq \sum_{k \geq 1} \mu(E \Delta T^{-k}E)$$

$$\leq \sum_{k \geq 1} \sum_{l=1}^k \mu(T^{-l+1}E \Delta T^{-l}E)$$

$$\stackrel{\text{mpt}}{=} \sum_{k \geq 1} \underbrace{\sum_{l=1}^k \mu(E \Delta T^{-1}E)}_{=0} = 0$$

If μ ergodic, then $\mu(E_0) = 0$ or $\mu(E_0^c) = 0$. Since $\mu(E \Delta E_0) = 0$, it follows $\mu(E) = 0$ or $\mu(E^c) = 0$.

"(b) \Rightarrow (c)": Let f be measurable s.t. $f = f \circ T$ a.e. For $t \in \mathbb{R}$,
 $E := \{f > t\}$ is measurable, and we have

$$T^{-1}E = \{x \in X \mid f(Tx) > t\} \stackrel{\text{a.e.}}{=} \{x \in X \mid f(x) > t\} = E,$$

i.e. $\mu(E \Delta T^{-1}E) = 0$. By assumption $\mu(\{f > t\}) = 0$ or $\mu(\{f \leq t\}) = 0$. Assume $f \neq \text{const}$ a.e. Then $\exists t^*$, s.t. $\mu(\{f > t^*\}) > 0$ and $\mu(\{f \leq t^*\}) > 0$, in contradiction with what we just obtained.

"(c) \Rightarrow (a)": Let $E \in \mathcal{B}$ be an invariant set. Then $f = \chi_E$ satisfies

$$f \circ T = \chi_E \circ T = \chi_{T^{-1}E} = \chi_E = f,$$

hence $f = \text{const}$ a.e. Thus $\mu(E) = 0$ or $\mu(E^c) = 0$, and ergodicity follows.

1.12 ERGODICITY: INVARIANT UNDER MEASURE-THEORETIC ISOMORPHISM

Prop: Assume that the mpts $(X_i, \mathcal{B}_i, \mu_i, T_i)$, $i=1,2$, are meas. th. isomorph. Then, if one is ergodic, so is the other.

Proof: Suppose the mpt #1 is ergodic.

Let $f: X_2 \rightarrow \mathbb{R}$ be measurable and s.t. $f \circ T_2 = f$ a.e.

Then $f \circ T_2 \circ \pi = f \circ \pi$ a.e.

" "
 $f \circ \pi \circ T_1$, so by Thm 1.11: $f \circ \pi = \text{const}$ a.e.

Since π is an essential bijection, $f = \text{const}$ a.e.

Thm 1.11
 \implies mpt #2 is ergodic



1.13 INDEPENDENCE AND MIXING

- (X, \mathcal{B}, μ) prob. space
- $E \in \mathcal{B}$: events
- $\mathbb{P}[X \in E] = \mu(E)$: probability law

$E, F \in \mathcal{B}$ are independent, if

$$\mu(E) = \underbrace{\mu(E|F)} := \frac{\mu(E \cap F)}{\mu(F)} \Leftrightarrow \mu(E \cap F) = \mu(E)\mu(F)$$

conditional prob.

Def: A ppt (X, \mathcal{B}, μ, T) is called (strongly) mixing if for every $E, F \in \mathcal{B}$ we have $\mu(E \cap T^{-k}F) \xrightarrow{k \rightarrow \infty} \mu(E)\mu(F)$

" $T^{-k}(F)$ is asymptotically independent of E " or

" $T^k(x)$ loses dependence on x for large k "

Exercise: Mixing is an invariant under measure-theoretic isom.

Prop A: Strong mixing implies ergodicity.

Proof: $E \in \mathcal{B}$ invariant: $E = T^{-k}(E) \forall k \geq 0$

$$\Rightarrow \mu(E) = \mu(E \cap T^{-k}E) \xrightarrow{k \rightarrow \infty} \mu(E)^2 \Rightarrow \mu(E) \in \{0, 1\} \stackrel{\mu(X)}{\parallel}$$

Prop B: Strong mixing $\Leftrightarrow \forall f, g \in L^2(X, \mu)$:

$$\int f \cdot g \circ T^k d\mu \xrightarrow{k \rightarrow \infty} \int f d\mu \cdot \int g d\mu$$

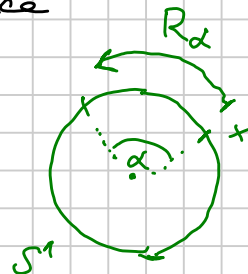
Proof: " \Leftarrow ": Take $f = \chi_E, g = \chi_F$

" \Rightarrow ": [Sa, Prop 1.3]

I.3 EXAMPLES

1.14 CIRCLE ROTATION

$\alpha \in \mathbb{R} : R_\alpha : S^1 \rightarrow S^1, \quad S^1: \text{circle with unit circumference}$
 $x \mapsto x + \alpha \pmod{1}$



\mathcal{B} : Borel σ -alg, $\mu = m$: Lebesgue measure

- Prop: (a) R_α is measure preserving
(b) R_α is ergodic iff $\alpha \notin \mathbb{Q}$
(c) R_α is never mixing

Proof: (a) \rightarrow Exercise

(b) $\bullet \alpha = \frac{p}{q} \in \mathbb{Q} \quad (p, q \in \mathbb{N})$

$\Rightarrow R_\alpha^q = \text{id}$

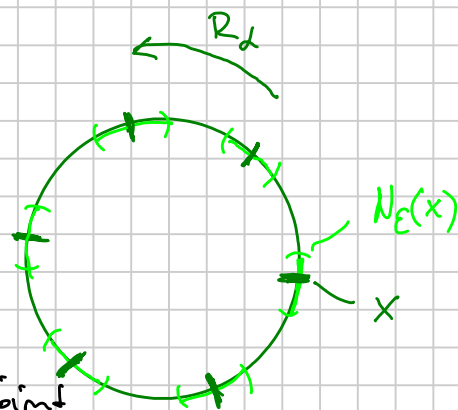
Pick $0 < \varepsilon < 1/2q$

$\Rightarrow N_\varepsilon(x), N_\varepsilon(x + \frac{1}{q}), \dots, N_\varepsilon(x + \frac{q-1}{q})$ disjoint

\uparrow " ε -neighborhood of x "

$\Rightarrow E = \bigcup_{i=0}^{q-1} N_\varepsilon(x + i \frac{1}{q})$ is invariant and $0 < m(E) < 1$

$\Rightarrow R_\alpha$ not ergodic



$\bullet \alpha \notin \mathbb{Q}$ Let $f = \chi_E, E$ invariant $\Rightarrow f \circ R_\alpha = f$

Fourier series: $f(x) = \sum_{n \in \mathbb{Z}} c_n e^{2\pi i n x}$ (in L^2)

$\rightarrow \sum_{n \in \mathbb{Z}} c_n e^{2\pi i n x} = f(x) = f(R_\alpha x) = \sum_{n \in \mathbb{Z}} c_n e^{2\pi i n (x + \alpha)}$

$\Leftrightarrow c_n = 0 \vee e^{2\pi i n \alpha} = 1 \quad \forall n \neq 0$

$\Leftrightarrow c_n = 0 \quad \forall n \neq 0 \Rightarrow f = \text{const} \Rightarrow m(E) \in \{0, 1\} \Rightarrow R_\alpha \text{ ergodic}$

(c) Set $E = (0, \varepsilon) = F$, $0 < \varepsilon < 1$

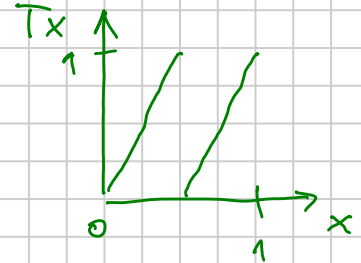
Dirichlet thm: $\forall \varepsilon \in \mathbb{R} \exists m_{\frac{\varepsilon}{2}} \xrightarrow{z \rightarrow \infty} \infty$ s.t. $m_{\frac{\varepsilon}{2}} d \bmod 1 \xrightarrow{(z \rightarrow \infty)} 0$

$$\Rightarrow m(E \cap R_{\frac{\varepsilon}{2}}^{-m_{\frac{\varepsilon}{2}}} F) \xrightarrow{z \rightarrow \infty} m(E) = \varepsilon \neq \varepsilon^2 = m(E)m(F)$$

$\Rightarrow \mathbb{R}_d$ not mixing

1.15 ANGLE DOUBLING

$$X = S^1, \mu = m, T(x) = 2x \bmod 1$$



Prop: T is measure preserving and strongly mixing

Proof: Exercise.

1.16 BERNOULLI SCHEMES

- S — alphabet (finite set)
- $X = S^{\mathbb{N}} \ni \{x_i\} = (x_0, x_1, x_2, \dots)$
- $d(\{x_i\}, \{y_i\}) = 2^{-\min\{i \mid x_i \neq y_i\}}$ — metric

↪ topology generated by cylinders

$$[a_0, a_1, \dots, a_{m-1}] = \{ \{x_i\} \mid x_i = a_i \ 0 \leq i \leq m-1 \}$$

Cylinders generate the Borel σ -alg w.r.t. metric d .

- $T: X \rightarrow X, (x_0, x_1, x_2, \dots) \mapsto (x_1, x_2, \dots)$ — left shift
- $p = (p_a)_{a \in S} \in \mathbb{R}^{|S|}$ — probability vector (i.e. $p \geq 0$
 $\sum_a p_a = 1$)

Def: The Bernoulli measure corresponding to p on \mathcal{B} (Borel σ -alg. wrt d) is the unique measure μ , s.t.

$$\mu[a_0, \dots, a_{m-1}] = \prod_{i=0}^{m-1} p_{a_i}$$

Prop A: The $(\frac{1}{2}, \frac{1}{2})$ -Bernoulli shift (i.e. $X = \{0, 1\}$, $p = (\frac{1}{2}, \frac{1}{2})$, μ the corresp. Bernoulli measure, T the left shift) is measure theoretically isomorphic to the angle doubling map.

Proof: Let $\pi(x_0, x_1, \dots) = \sum_{i=0}^{\infty} 2^{-(i+1)} x_i$

$\pi: X^{\mathbb{N}} \rightarrow S^1$ is bijective, where

$$X^{\mathbb{N}} = \{ \{x_i\} \in X^{\mathbb{N}} \mid \nexists m \in \mathbb{N} \text{ s.t. } x_l = 1 \ \forall l \geq m \}$$

and $\mu(X^{\mathbb{N}}) = \mu(X) = 1$ (exercise)

Clearly $\pi \circ T = \tilde{T} \circ \pi$ with $\tilde{T}: S^1 \rightarrow S^1, x \mapsto 2x \pmod{1}$

Check that $m(\pi[a_0, \dots, a_{m-1}]) = 2^{-m} = \mu[a_0, \dots, a_{m-1}]$
 Mon. class thm.

$\implies \pi$ preserves the measure $\xrightarrow{\text{Def. 1.7C}}$ claim

\uparrow cylinders span $\mathcal{B}(X)$

\uparrow dyadic intervals span $\mathcal{B}(S^1)$ ▀

Prop B: Every Bernoulli scheme is mixing, hence ergodic.

Proof: The proof goes as in the case of the angle doubling map.

For the $(\frac{1}{2}, \frac{1}{2})$ -Bernoulli shift we could use Prop 1.15 + Prop A (here) + the invariance of mixing under isomorphism to get the claim. ▀

$\{x_k\} \in X$

sequence of events Y_0, Y_1, \dots

$Y_i \in \{0, 1\}$
heads tails

$\mu[a_0, \dots, a_{m-1}] = 2^{-m}$

$\mathbb{P}[Y_0 = a_0, \dots, Y_{m-1} = a_{m-1}] = 2^{-m}$

Invariance: $\mu = \mu \circ T^{-1}$

$\Rightarrow \mu[a_0] = \mu[* , a_0]$

$\mathbb{P}[Y_0 = a_0] = \mathbb{P}[Y_0 = *, Y_1 = a_0]$
 $= \mathbb{P}[Y_1 = a_0]$

independence of the Y_i

1.17 SUBSHIFT OF FINITE TYPE

S - alphabet (finite),

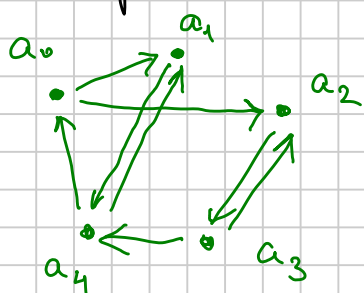
$A = \{t_{ij}\}_{i,j \in S} \in \{0, 1\}^{|S| \times |S|}$

Def: The subshift of finite type (sft) with alphabet S and transition matrix A is:

- $\Sigma_A^+ = \{ \{x_k\} \in S^{\mathbb{N}} \mid t_{x_i x_{i+1}} = 1 \ \forall i \geq 0 \}$, and
- metric $d(\{x_k\}, \{y_k\}) = 2^{-\min\{k \mid x_k \neq y_k\}}$, and
- the action $T(x_0, x_1, x_2, \dots) = (x_1, x_2, \dots)$, the left shift

The sft is a compact metric space, and $T: \Sigma_A^+ \rightarrow \Sigma_A^+$ is continuous

Idea:



$a_i \in S, i = 1, \dots, |S|$

Edge from a_i to a_j iff $t_{a_i a_j} = 1$

$\Rightarrow \Sigma_A^+$: the set of all possible paths on this graph

Further terminology:

• Stochastic matrix: $P = (P_{ab})_{a,b \in S} \in [0, 1]^{|S| \times |S|}, P \mathbb{1} = \mathbb{1} = (1, \dots, 1)^T$

compatible with A: $t_{ab} = 0 \Rightarrow P_{ab} = 0$

• Probability vector: $P = (p_a)_{a \in S} \in [0, 1]^{|S|}$ s.t. $\sum p_a = 1$ 15

• Stationary probability vector: p s.t. $p^T P = p^T$

1.18 MARKOV SHIFT

Def: Given: prob. vector p , stoch. matrix P compatible with A .
The Markov measure on Σ_A^+ (with the Borel σ -alg $\mathcal{B}(\Sigma_A^+)$) is defined through

$$\mu[a_0, a_1, \dots, a_{m-1}] = p_{a_0} \cdot p_{a_0 a_1} \cdot \dots \cdot p_{a_{m-2} a_{m-1}}$$

(Remark: P stochastic $\Rightarrow \mu$ is well-defined Borel prob. measure)

Then, $(\Sigma_A^+, \mathcal{B}(\Sigma_A^+), \mu, T)$ is called the (p, P) -Markov shift.
 \uparrow left shift

Prop: (a) p stationary wrt $P \Leftrightarrow (p, P)$ -Markov shift is a ppt

(b) Every stoch. matrix has a stat. prob. vector

Proof: (a) Sufficient to show for cylinders: $\mu[\underline{b}] = \mu(T^{-1}[\underline{b}]) = \mu[\ast, \underline{b}]$:

$$\sum_a p_a p_{a b_0} \cdot p_{b_0 b_1} \cdot \dots \cdot p_{b_{m-2} b_{m-1}} = p_{b_0} p_{b_0 b_1} \cdot \dots \cdot p_{b_{m-2} b_{m-1}}$$

no division by 0 due to compatibility $\Leftrightarrow \sum_a p_a p_{a b_0} = p_{b_0}$

(b) Let $\Delta = \{x \in \mathbb{R}^{|S|} \mid x \geq 0, \sum x_i = 1\}$, and

$$R: \Delta \rightarrow \mathbb{R}^{|S|}, x \mapsto (x^T P)^T = P^T x$$

Check that $R(\Delta) \subseteq \Delta$ due to stochasticity of P

Brouwer's fixed point theorem (continuous mapping on a convex, compact set has a fixed point) $\Rightarrow R$ has a fixed point

Ergodicity & mixing of Markov shifts \rightarrow later

I.4 BASIC CONSTRUCTIONS

1.19 PRODUCTS

$(X_i, \mathcal{B}_i, \mu_i)$, $i=1,2$, measure spaces

$(X_1 \times X_2, \mathcal{B}_1 \otimes \mathcal{B}_2, \mu_1 \times \mu_2)$ — product (measure) space

• $\mathcal{B}_1 \otimes \mathcal{B}_2 = \sigma(\{B_1 \times B_2 \mid B_i \in \mathcal{B}_i\})$

• $(\mu_1 \times \mu_2)(B_1 \times B_2) := \mu_1(B_1)\mu_2(B_2)$ (unique!)

Def: The product of two mpts $(X_i, \mathcal{B}_i, \mu_i, T_i)$, $i=1,2$, is the mpt $(X_1 \times X_2, \mathcal{B}_1 \otimes \mathcal{B}_2, \mu_1 \times \mu_2, T_1 \times T_2)$ with

$$(T_1 \times T_2)(x_1, x_2) = (T_1 x_1, T_2 x_2)$$

Prop: The product of two ergodic mpts is not necessarily ergodic.
The product of two mixing mpts is always mixing.

Proof: The product of two (identical) ergodic circle rotations:

$$T: S^1 \times S^1 \rightarrow S^1 \times S^1, (x_1, x_2) \mapsto (x_1 + d, x_2 + d) \pmod{1}$$

" $\mathbb{R}_d \times \mathbb{R}_d$

is not ergodic: $f(x,y) := x-y \pmod{1}$ is a non-constant invariant function (cf. Thm. 1.11)

Mixing: see [Sa, Prop. 1.9]

1.20 SKEW-PRODUCTS

Example: $(\Sigma, \mathcal{B}, \mu, \sigma)$: $(\frac{1}{2}, \frac{1}{2})$ -Bernoulli shift ($\Sigma = \{0,1\}^{\mathbb{N}}$)

$$f: \Sigma \rightarrow \mathbb{Z}, f(\{x_k\}) = (-1)^{x_0}$$

$$\text{Let } T: \Sigma \times \mathbb{Z} \rightarrow \Sigma \times \mathbb{Z}, T(\{x_k\}, l) = (\sigma\{x_k\}, l + f(\{x_k\}))$$

T is mpt wrt $\mu \times m_{\mathbb{Z}}$ ($m_{\mathbb{Z}}$: counting measure), and 17

$$T^m(\{x_i\}, l) = (\sigma^m\{x_i\}, l + Y_0 + Y_1 + \dots + Y_{m-1}), \quad Y_i := (-1)^{x_i}$$

"random walk on \mathbb{Z} driven by the noise process $(X, \mathcal{B}, \mu, \sigma)$ "

Def: Let $(\Omega, \mathcal{F}, \rho, \sigma)$ be a mpt, (X, \mathcal{B}, μ) a prob. space, and $\{T_\omega\}_{\omega \in \Omega}$ a family of mpts on (X, \mathcal{B}, μ) , s.t. $(\omega, x) \mapsto T_\omega x$ is measurable with respect to $\rho \times \mu$.

We define the **skew-product** as $(\Omega \times X, \mathcal{F} \otimes \mathcal{B}, \rho \times \mu, \tau)$ with $\tau(\omega, x) = (\sigma\omega, T_\omega x)$

τ preserves $\rho \times \mu \rightsquigarrow$ exercise (use Fubini's theorem)

$$\tau^m(\omega, x) = (\sigma^m\omega, T_{\sigma^{m-1}\omega} \circ T_{\sigma^{m-2}\omega} \circ \dots \circ T_{\sigma\omega} \circ T_\omega x)$$

\rightsquigarrow skew-products model forced/driven/time-dependent/random dyn. systems

The skew-product is sometimes called a **random dynamical system**

$(\Omega, \mathcal{F}, \rho, \sigma)$ is also called **base** or **driving system**

Example: T_ω : change in the amount of ice in a glacier

ω : season ($\sigma^m\omega$ models planetary dynamics \rightsquigarrow amount of incoming solar radiation)

1.21 FACTORS AND EXTENSIONS

Def: A mpt (X, \mathcal{B}, μ, T) is called a **factor** of a mpt (Y, \mathcal{C}, ν, S) if $\exists X' \subseteq X, \exists Y' \subseteq Y$, both of full measure, s.t.

- $T(X') \subseteq X'$, $S(Y') \subseteq Y'$
- there is a $\pi: Y' \rightarrow X'$ measurable and onto (surjective), with $\nu \circ \pi^{-1} = \mu$ and $\pi \circ S = T \circ \pi$ on Y'

π is called the **factor map**. (Y, \mathcal{C}, ν, S) is called an **extension** of (X, \mathcal{B}, μ, T) .

Examples

- 1) Meas.-theoretically isomorphic systems are factors and extensions of each other
- 2) A skew-product $\tau: \Omega \times X \rightarrow \Omega \times X$ is an extension of its base $\sigma: X \rightarrow X$, if μ is a probability measure.

3) Assume, we can "measure" the state $x \in X$ of some mpst only through the **observables** $f: X \rightarrow \mathbb{R}$ (measurable). The dynamical information we obtain is encoded in the sub- σ -alg

$$\mathcal{B} := \sigma \left(f \circ T^m \mid m \geq 0 \right) = \sigma \left\{ T^{-m} \circ f^{-1}(A) \mid A \in \mathcal{B}(\mathbb{R}) \right\}$$

(E.g. $f = \chi_A$ for some $A \in \mathcal{C}$, and then we observe as trajectories sequences in $\{0,1\}^{\mathbb{N}}$)

It can happen that (X, \mathcal{B}, μ, T) is mixing, while (X, \mathcal{C}, μ, T) isn't.