

STOCHASTIC DYNAMICS

LECTURE COURSE
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1. Markov chains: construction and elementary properties

Definition 1.1. Let (S, \mathcal{S}) be a measurable space. A function

$$p : S \times \mathcal{S} \rightarrow [0, 1]$$

is called *transition probability*, if

- (a) for each $x \in S$ $p(x, \cdot)$ is a probability measure on (S, \mathcal{S}) ;
- (b) for each $A \in \mathcal{S}$ $p(\cdot, A)$ is \mathcal{S} -measurable.

Remark 1.2. Let p be a transition probability on a measurable space (S, \mathcal{S}) .

- i) If $f : S \rightarrow \mathbb{R}$ is \mathcal{S} - \mathcal{B}^1 -measurable and bounded, then also $g := \int_S f(x) p(\cdot, dx)$ is.
- ii) If μ is a probability measure on (S, \mathcal{S}) , then also $\nu := \int_S p(x, \cdot) \mu(dx)$ is.

PROOF. i) By 1.1(b) the statement holds for indicator functions $f = 1_A$ of measurable sets, hence also (linearity of the integral) for step functions.

If $f \geq 0$, there exist approximating step functions $0 \leq f_n \nearrow f$; here $g_n := \int f_n(x) p(\cdot, dx)$ is measurable (since f_n is a step function) and (by the bound of f) bounded; on the other hand (monotone convergence) $g_n \nearrow \int f(x) p(\cdot, dx) \equiv g$, such that g is measurable as pointwise limit of measurable functions; g is bounded with the same bound as f .

More generally, decompose $f = f^+ - f^-$ mit $f^+, f^- \geq 0$; by what we know, $g^\pm := \int f^\pm(x) p(\cdot, dx)$ measurable and bounded, hence also $g = g^+ - g^-$.

ii) For a sequence $(A_n)_{n \in \mathbb{N}}$ of pairwise disjoint $A_n \in \mathcal{S}$ we have

$$\begin{aligned} \nu \left(\bigcup_n A_n \right) &\equiv \int_S p \left(x, \bigcup_n A_n \right) \mu(dx) \stackrel{1.1(a)}{=} \int_S \sum_n p(x, A_n) \mu(dx) \\ &\stackrel{\text{mon. conv.}}{=} \sum_n \int_S p(x, A_n) \mu(dx) \equiv \sum_n \nu(A_n) ; \end{aligned}$$

moreover $\nu(S) \equiv \int_S p(x, S) \mu(dx) \stackrel{1.1(a)}{=} \int_S \mu(dx) = 1$ (μ probability measure). \square

Definition 1.3. Let S be a Polish space, and $(p_n)_{n \in \mathbb{N}}$ a sequence of transition probabilities, μ a probability measure. Then let $P_0 := \mu$ and

$$P_n(B_0 \times \cdots \times B_n) := \int_{B_0 \times \cdots \times B_n} p_n(x_{n-1}, dx_n) p_{n-1}(x_{n-2}, dx_{n-1}) \cdots p_1(x_0, dx_1) \mu(dx_0)$$

for $n \in \mathbb{N}$ and $B_i \in \mathcal{S} \equiv \mathcal{B}(S)$.

By remark 1.2 we obtain recursively that P_n is well defined on the semiring¹

$$\mathcal{R}_n := \{B_0 \times \cdots \times B_n : B_i \in \mathcal{S}\}.$$

Moreover, P_n can be extended to a measure on the σ -algebra generated by the ring $r(\mathcal{R}_n)$

$$\sigma(r(\mathcal{R}_n)) = \underbrace{\mathcal{S} \otimes \cdots \otimes \mathcal{S}}_{(n+1)\text{times}} \equiv \mathcal{S}^{n+1} \quad \left(\overset{S \text{ Polish}}{=} \mathcal{B}(S^{n+1}) \right).$$

PROOF. P_n induces a pre-measure on the ring $r(\mathcal{R}_n)$. By Caratheodory, to extend it to $\sigma(r(\mathcal{R}_n))$ we have to prove that P_n is σ -additive on the ring. Since P_n is a finite pre-measure, σ -additivity is equivalent to the property of *continuity from above* that will be shown in what follows; by recursion and remark 1.2 it is enough to consider the case $n = 1$:

Let $(A_k)_{k \in \mathbb{N}}$ be a sequence in $r(\mathcal{R}_1)$ with $A_k \searrow \emptyset$. We have to prove: $P_1(A_k) \xrightarrow{k \rightarrow \infty} 0$. Denoting the intersection by $A \in r(\mathcal{R}_1)$ near $x \in S$ with

$$A_x := \{y \in S : (x, y) \in A\},$$

we obtain by $A_k \searrow \emptyset$ for all $x \in S$:

$$(A_k)_x \searrow \emptyset \quad (k \rightarrow \infty).$$

By *continuity from above* of the measure $p_1(x, \cdot)$ we get

$$p_1(x, (A_k)_x) \xrightarrow{k \rightarrow \infty} 0 \quad (x \in S)$$

and thus by dominated convergence:

$$P_1(A_k) = \int_S p_1(x, (A_k)_x) \mu(dx) \xrightarrow{k \rightarrow \infty} 0.$$

□

Our next aim is the construction of a Markov chain on $S^{\mathbb{N}_0}$ with transition probabilities $(p_n)_{n \in \mathbb{N}}$ and initial distribution μ . For this, let S be Polish, endowed with the Borel σ -algebra $\mathcal{B}(S) =: \mathcal{S}$. We verify the consistency condition by Kolmogorov for $(P_n)_{n \in \mathbb{N}_0}$. For this purpose we define:

For $F, G \subset \mathbb{N}_0$ with $F \subset G$ let

$$\begin{aligned} \pi_{G,F} : S^G &\longrightarrow S^F \\ (x_i)_{i \in G} &\longmapsto (x_i)_{i \in F} \end{aligned}$$

¹A collection of sets is called \mathcal{P} *semiring*, if (cf. Halmos [HM 74, S.22]):

- for $E \in \mathcal{P}$ and $F \in \mathcal{P}$ also $E \cap F \in \mathcal{P}$, and
- for $E \in \mathcal{P}$ and $F \in \mathcal{P}$ mit $E \subset F$ there exist pairwise disjoint sets $C_1, \dots, C_n \in \mathcal{P}$, such that

$$F \setminus E = \cup_{i=1}^n C_i.$$

the projection on the smaller index set and herewith $\pi_F := \pi_{\mathbb{N}_0, F}$; correspondingly let $m, n \in \mathbb{N}_0$ with $m \leq n$

$$\begin{aligned} \pi_{n,m} : S^{n+1} &\longrightarrow S^{m+1} \\ (x_0, \dots, x_n) &\longmapsto (x_0, \dots, x_m) \end{aligned}$$

and for $m \in \mathbb{N}_0$

$$\begin{aligned} \pi_m : S^{\mathbb{N}_0} &\longrightarrow S^{m+1} \\ (x_i)_{i \in \mathbb{N}_0} &\longmapsto (x_0, \dots, x_m). \end{aligned}$$

These projections are measurable for the respective product σ -algebras. The uniqueness theorem for measures (applied to \cap -stable generators of σ -algebras consisting of cylinder sets) yields for $m, n \in \mathbb{N}_0$ with $m \leq n$ the equality

$$P_n \circ \pi_{n,m}^{-1} = P_m$$

of measures on \mathcal{S}^{m+1} ; analogously for finite $F, G \subset \mathbb{N}_0$ with $F \subset G$ also

$$P_G \circ \pi_{G,F}^{-1} = P_{\max F} \circ \pi_{\{0, \dots, \max F\}, F}^{-1} =: P_F;$$

this consistency property implies that $(P_F)_{F \subset \mathbb{N}_0 \text{ finite}}$ defines a pre-measure on $(S^{\mathbb{N}_0}, \mathcal{B}(S)^{\mathbb{N}_0})$. By Kolmogorov's consistency theorem it is even σ -additive. Therefore there is a unique probability measure P_μ on $(S^{\mathbb{N}_0}, \mathcal{B}(S)^{\mathbb{N}_0})$ with

$$P_\mu \circ \pi_n^{-1} = P_n \quad (n \in \mathbb{N}_0). \quad (1)$$

Satz 1.4 (canonical Markov chain). *Let S be a Polish space with transition probabilities $(p_n)_{n \in \mathbb{N}}$ and probability measure μ ; let P_μ be the thus induced probability measure on $S^{\mathbb{N}_0}$. Then*

$$X_n := \pi_{\{n\}} \equiv \pi_{\mathbb{N}_0, \{n\}} \quad (n \in \mathbb{N}_0)$$

is a Markov chain on

$$(\Omega, \mathcal{F}, \mathbb{P}, (\mathcal{F}_n)_{n \in \mathbb{N}_0}) := \left(S^{\mathbb{N}_0}, \mathcal{S}^{\mathbb{N}_0}, P_\mu, (\sigma(\pi_n))_{n \in \mathbb{N}_0} \right),$$

i.e. we have

i) X_n is \mathcal{F}_n -measurable, and

ii) for all $n \in \mathbb{N}_0$ and $B \in \mathcal{S}$ we have:

$$\mathbb{P}(X_{n+1} \in B \mid \mathcal{F}_n) = \mathbb{P}(X_{n+1} \in B \mid X_n) = p_{n+1}(X_n, B).$$

PROOF. i) X_n is measurable with respect to $\sigma(X_n) \subset \sigma(X_0, \dots, X_n) \equiv \sigma(\pi_n) \equiv \mathcal{F}_n$.

ii) We have to show:

$$\int_A 1_{\{X_{n+1} \in B\}} dP_\mu = \int_A p_{n+1}(X_n, B) dP_\mu \quad (A \in \mathcal{F}_n)$$

(then also $\mathbb{P}(X_{n+1} \in B | X_n) = p_{n+1}(X_n, B)$, by measurability with respect to $\sigma(X_n) \subset \mathcal{F}_n$). Since $\pi_n^{-1}(\mathcal{B}_n)$ is a \cap -stable generator of \mathcal{F}_n , it is enough, to prove the above equation for

$$A = \pi_n^{-1}(B_0 \times \dots \times B_n) \equiv \{X_0 \in B_0, \dots, X_n \in B_n\}$$

with $B_0, \dots, B_n \in \mathcal{S}$, namely:

$$\begin{aligned} \int_A 1_{\{X_{n+1} \in B\}} dP_\mu &= P_\mu\{X_0 \in B_0, \dots, X_n \in B_n, X_{n+1} \in B\} \\ &\stackrel{(1)}{=} P_{n+1}(B_0 \times \dots \times B_n \times B) \\ &\stackrel{1.3}{=} \int_{B_0 \times \dots \times B_n} p_{n+1}(x_n, B) P_n(dx_0, \dots, dx_n) \\ &\stackrel{\text{transf. thm}}{=} \int_A p_{n+1}(X_n, B) dP_\mu. \end{aligned}$$

□

Definition 1.5. Let (S, \mathcal{S}) be a measurable space. Then on the path space $\Omega \equiv S^{\mathbb{N}_0}$ the family $\theta \equiv (\theta_n)_{n \in \mathbb{N}_0}$ of (*canonical*) *shifts* $\theta_n : \Omega \longrightarrow \Omega$ ($n \in \mathbb{N}_0$) is defined by

$$\theta_n(\omega) := (m \mapsto \omega(m+n)).$$

Each θ_n is measurable with respect to $\mathcal{F} \equiv \mathcal{S}^{\mathbb{N}_0}$.

We next prove the Markov property (with deterministic times) and then the strong Markov property (with stopping times). For this we denote \mathbb{E}_μ resp. $\mathbb{E}_x \equiv \mathbb{E}_{\delta_x}$ the conditional expectations with respect to P_μ resp. P_{δ_x} on Ω with underlying transition probabilities $(p_n)_{n \in \mathbb{N}}$. For simplicity we assume that the Markov chain is time homogeneous:

Definition 1.6. In the situation of Proposition 1.4 the Markov chain is called *X time homogeneous*, if for all $n \in \mathbb{N}$ we have $p_n = p_1 (= p)$.

Theorem 1.7 (Markov property). *In the situation of 1.4 let the Markov chain X be time homogeneous; let Y be a bounded, \mathcal{F} -measurable random variable on Ω . Then*

$$\mathbb{E}_\mu(Y \circ \theta_n | \mathcal{F}_n) = \mathbb{E}_{X_n}(Y) \equiv \mathbb{E}_x(Y) \Big|_{x=X_n} \quad (n \in \mathbb{N}_0).$$

PROOF. Note first that $\mathbb{E}_{X_n}(Y)$ is indeed measurable with respect to \mathcal{F}_n ; this follows from adaptedness of X and measurability of $x \mapsto \mathbb{E}_x(Y)$ [this one by definition and recursive application of 1.2 i) clear for indicator functions $Y = 1_{\pi_n^{-1}[B_0 \times \dots \times B_n]}$ for $B_i \in \mathcal{S}$; the

general property follows from the monotone class theorem, since by monotone convergence $\{Y : x \mapsto \mathbb{E}_x(Y) \text{ measurable}\}$ is closed for monotone operations]. So we have to prove the claimed equation. By monotone class arguments, it is enough to argue for the case Y of the form $\prod_{k=0}^m g_k(X_k)$ with bounded \mathcal{S} -measurable random variables g_0, \dots, g_m .

1) We first consider sets \mathcal{F}_n of the form $A := \pi_n^{-1}[A_0 \times \dots \times A_n]$ with $A_0, \dots, A_n \in \mathcal{S}$; we have

$$\begin{aligned} \mathbb{E}_\mu(Y \circ \theta_n \cdot 1_A) &\equiv \mathbb{E}_\mu \left(\prod_{k=0}^m g_k(X_{n+k}) \cdot 1_A \right) \\ &\stackrel{(1),1.3}{=} \int_{A_0} \mu(dx_0) \int_{A_1} p(x_0, dx_1) \cdots \int_{A_n} p(x_{n-1}, dx_n) \times \\ &\quad \times \int_S g_0(x_{n+1}) p(x_n, dx_{n+1}) \cdots \int_S g_m(x_{n+m}) p(x_{n+m-1}, dx_{n+m}) \\ &\stackrel{\text{transf. thm}}{=} \mathbb{E}_\mu \left(\mathbb{E}_{X_n} \left(\prod_{k=0}^m g_k(X_k) \right) \cdot 1_A \right) \\ &\equiv \mathbb{E}_\mu(\mathbb{E}_{X_n}(Y) \cdot 1_A), \end{aligned}$$

hence the claim follows for all $A \in \mathcal{F}_n$ of the special considered form.

2) Let now $\mathcal{L} := \{A \in \mathcal{F}_n : \text{claim from 1) valid for } A\}$. According to 1) $\pi_n^{-1}(\mathcal{R}_n) \subset \mathcal{L}$; since $\pi_n^{-1}(\mathcal{R}_n)$ is \cap -stable, Dynkin's lemma yields $\mathcal{F}_n = \sigma(\pi_n^{-1}(\mathcal{R}_n)) \subset \mathcal{L}$. \square

Our next goal is to extend the Markov property to stopping times.

Definition 1.8. Let $\Omega, \mathcal{F}, (\mathcal{F}_n)_{n \in \mathbb{N}_0}$ be a filtered measure space; $N : \Omega \rightarrow \mathbb{N}_0 \cup \{\infty\}$ is called $(\mathcal{F}_n)_{n \in \mathbb{N}_0}$ -stopping time, if $\{N \leq n\} \in \mathcal{F}_n$ for all \mathbb{N}_0 . Equivalently, $\{N = n\} \in \mathcal{F}_n$ for all $n \in \mathbb{N}_0$.

To an $(\mathcal{F}_n)_n$ -stopping time N we associate the σ -algebra

$$\mathcal{F}_N := \left\{ A \in \mathcal{F} : A \cap \left\{ N \stackrel{(\equiv)}{\leq} n \right\} \in \mathcal{F}_n \text{ for all } n \in \mathbb{N}_0 \right\};$$

it is called N -past or σ -algebra of events before N .

In the situation from 1.4 and 1.5 we formally enlarge Ω by $\Delta \notin \Omega$, increase \mathcal{F} by $\{\Delta\}$ and define a $(\mathcal{F}_n)_{n \in \mathbb{N}_0}$ -stopping time N by

$$\theta_N(\omega) := \begin{cases} \theta_{N(\omega)}(\omega) & , N(\omega) < \infty \\ \Delta & , N(\omega) = \infty . \end{cases}$$

For a random variable Y on Ω let $Y(\Delta) := 0$.

Theorem 1.9 (strong Markov property). *In the situation from 1.4 assume the Markov chain X is time homogeneous; let θ be the shift from 1.5 and N a $(\mathcal{F}_n)_n$ -stopping time. If $(Y_n)_{n \in \mathbb{N}_0}$ is a family of \mathcal{F} -measurable and (uniformly in (n, ω)) bounded random variables, we have*

$$\mathbb{E}_\mu(Y_N \circ \theta_N | \mathcal{F}_N) = \mathbb{E}_{X_N}(Y_N) \quad \text{on } \{N < \infty\};$$

In particular, for a \mathcal{F} -measurable bounded random variable Y :

$$\mathbb{E}_\mu(Y \circ \theta_N | \mathcal{F}_N) = \mathbb{E}_{X_N}(Y) \quad \text{on } \{N < \infty\}.$$

PROOF. Note first that $\omega \mapsto \mathbb{E}_{X_{N(\omega)}(\omega)}(Y_{N(\omega)})$ is in fact \mathcal{F}_N -measurable, since it is the composition of the measurable maps $\omega \mapsto (\omega, N(\omega))$, $(\omega, n) \mapsto (X_n(\omega), n)$ and $(x, n) \mapsto \mathbb{E}_x(Y_n)$.

With $A \in \mathcal{F}_N$ we then have

$$\begin{aligned} \mathbb{E}_\mu(Y_N \circ \theta_N \cdot 1_{A \cap \{N < \infty\}}) &\stackrel{\text{dom. conv.}}{=} \sum_{n=0}^{\infty} \mathbb{E}_\mu(Y_n \circ \theta_n \cdot 1_{A \cap \{N=n\}}) \\ &\stackrel{\text{MP 1.7}}{=} \sum_{n=0}^{\infty} \mathbb{E}_\mu(\mathbb{E}_{X_n}(Y_n) \cdot 1_{A \cap \{N=n\}}) \\ &\stackrel{\text{dom. conv.}}{=} \mathbb{E}_\mu(\mathbb{E}_{X_N}(Y_N) \cdot 1_{A \cap \{N < \infty\}}). \end{aligned}$$

□

We next aim at investigating invariant measures of a Markov chain. Invariant measures are strongly correlated with return properties. We therefore assume as a further simplification that S is countable; for the representation with general Polish state space S see Meyn & Tweedie [M-T 93].

Let in the following

$$T_y := \inf\{n \in \mathbb{N} : X_n = y\} \quad (y \in S),$$

be the first hitting time of y and thus

$$\rho_{xy} := P_x(T_y < \infty) \quad (x, y \in S).$$

$y \in S$ is called $\left\{ \begin{array}{l} \text{recurrent} \\ \text{transient} \end{array} \right\}$, if $\left\{ \begin{array}{l} \rho_{yy} = 1 \\ \rho_{yy} < 1 \end{array} \right\}$. The number of visits in y ,

$$H_y := \sum_{n=1}^{\infty} 1_{\{X_n=y\}}$$

characterizes recurrence and transience of y in the following way:

Theorem 1.10 (transience and recurrence). *Let the Markov chain X from 1.4 be time homogeneous with countable state space S . Then for $y \in S$:*

$$\begin{aligned} y \text{ transient} &\implies \mathbb{E}_x(H_y) = \frac{\rho_{xy}}{1 - \rho_{yy}} < \infty \quad (\forall x \in S), \\ y \text{ recurrent} &\iff \mathbb{E}_y(H_y) = \infty. \end{aligned}$$

PROOF. For $k \in \mathbb{N}$ let T_y^k the time of the k -th visit in y . With this we have

$$P_x(T_y^k < \infty) = \rho_{xy} \cdot \rho_{yy}^{k-1} \quad (x \in S, k \in \mathbb{N}); \quad (\star)$$

for $k = 1$ this is just the definition of ρ_{xy} ; for $k > 1$ it follows inductively:

$$\begin{aligned} P_x(T_y^k < \infty) &= P_x\left(T_y^{k-1} < \infty, T_y \circ \theta_{T_y^{k-1}} < \infty\right) \\ &= \mathbb{E}_x\left(1_{\{T_y^{k-1} < \infty\}} \underbrace{\mathbb{E}_x\left(1_{\{T_y \circ \theta_{T_y^{k-1}} < \infty\}} \mid \mathcal{F}_{T_y^{k-1}}\right)}_{\text{str. MP 1.9 } \mathbb{E}_y(1_{\{T_y < \infty\}}) \equiv \rho_{yy}}\right) \\ &= \rho_{yy} \cdot P_x(T_y^{k-1} < \infty) \\ &\stackrel{\text{ind. hyp.}}{=} \rho_{xy} \rho_{yy}^{k-1}. \end{aligned}$$

Therefore

$$\mathbb{E}_x(H_y) = \sum_{n=1}^{\infty} P_x\left\{\underbrace{H_y \geq n}_{\{T_y^n < \infty\}}\right\} \stackrel{(\star)}{=} \rho_{xy} \sum_{n=1}^{\infty} \rho_{yy}^{n-1} = \frac{\rho_{xy}}{1 - \rho_{yy}};$$

the geometric series converges for $\rho_{yy} < 1$, and diverges iff $\rho_{yy} = 1$. \square

We next show that recurrence is contagious:

Theorem 1.11. *Let the Markov chain X from 1.4 be time homogeneous with countable S . If $x \in S$ is recurrent and $\rho_{xy} > 0$ with some $y \in S$, then also y is recurrent and we have $\rho_{yx} = 1$.*

PROOF. By recurrence of von x we have

$$\begin{aligned} 0 &= P_x(T_x = \infty) \geq P_x(T_y < \infty, T_x \circ \theta_{T_y} = \infty) \\ &= \mathbb{E}_x \left(1_{\{T_y < \infty\}} \underbrace{\mathbb{E}_x \left(1_{\{T_x \circ \theta_{T_y} = \infty\}} \mid \mathcal{F}_{T_y} \right)}_{\stackrel{\text{str. MP 1.9}}{=} \mathbb{E}_y(1_{\{T_x = \infty\}}) = (1 - \rho_{yx})} \right) \\ &= \rho_{xy} (1 - \rho_{yx}) ; \end{aligned}$$

Since by hypothesis $\rho_{xy} > 0$, we obtain: $\rho_{yx} = 1$.

With this the recurrence of y follows: by $\rho_{xy} > 0$ and $\rho_{yx} = 1$ there exist $k_1, k_2 \in \mathbb{N}$ with

$$P_x(X_{k_1} = y) > 0 \quad \text{and} \quad P_y(X_{k_2} = x) > 0 .$$

By Chapman-Kolmogorov for $n \in \mathbb{N}$ we have:

$$P_y(X_{n+k_1+k_2} = y) \geq P_y(X_{k_2} = x) P_x(X_n = x) P_x(X_{k_1} = y) ,$$

hence

$$\mathbb{E}_y(H_y) = \sum_{n=1}^{\infty} P_y(X_n = y) \geq \underbrace{P_y(X_{k_2} = x)}_{>0} \underbrace{\mathbb{E}_x(H_x)}_{\stackrel{1.10}{=} \infty} \underbrace{P_x(X_{k_1} = y)}_{>0} .$$

Hence also $\mathbb{E}_y(H_y) = \infty$ and y is recurrent by 1.10. \square

Thus the set of recurrent states decomposes into classes: For $x, y \in S$ let

$$x \sim y \quad :\iff \quad (x = y \text{ or } (\rho_{xy} > 0 \text{ and } \rho_{yx} > 0)) .$$

Theorem 1.12. *Let the Markov chain X from 1.4 be time homogeneous with countable S . Then the set of recurrent states $R := \{x \in S : \rho_{xx} = 1\}$ decomposes into a family $(R_i)_{i \in I}$ of pairwise disjoint classes, the equivalence classes of \sim .*

PROOF. We have to show that \sim is an equivalence relation: reflexivity and symmetry follow directly from the definition, so that only transitivity remains to prove:

If $x, y, z \in R$ are fixed, we have to show that with $x \sim y$ and $y \sim z$ also $x \sim z$ holds true. For this purpose we may wlog assume $x \neq y$ and $x \neq z$; by definition of \sim we have $\rho_{xy} > 0$ and $\rho_{yz} > 0$. Applying the strong Markov property as in the proofs of 1.10 and 1.11 we obtain:

$$\rho_{xz} \equiv P_x(T_z < \infty) \geq P_x(T_y < \infty, T_z \circ \theta_{T_y} < \infty) = \rho_{xy} \rho_{yz} > 0 ,$$

whence with 1.11 we get $(x \in R) : \rho_{zx} = 1 > 0$, in summary $x \sim z$. \square

2. Invariant measures and asymptotic behavior

We further consider the following situation: The countable space S is state space of a canonical time homogeneous Markov chain $(X_n)_{n \in \mathbb{N}_0}$ with space of trajectories $(\Omega, \mathcal{F}) := (S^{\mathbb{N}_0}, \mathcal{S}^{\mathbb{N}_0})$ and transition matrix p .

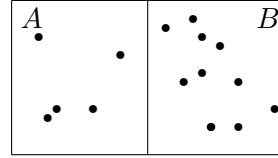
Definition 2.1. A measure μ on \mathcal{S} is called *stationary*, if for all $y \in S$ we have

$$\mu(y) = (\mu p)(y) \equiv \sum_{x \in S} \mu(x) p(x, y) < \infty .$$

A measure μ on \mathcal{S} is called *invariant*, if it is a stationary probability measure.

Example 2.2 (Ehrenfest model of diffusions). In a system consisting of the containers A and B we have a total of r molecules. Let X_n be the number of molecules in A at time $n \in \mathbb{N}_0$. This quantity takes its values in $S := \{0, 1, \dots, r\}$. By

$$p(k, m) := \begin{cases} \frac{r-k}{r} & , m = k + 1 \\ \frac{k}{r} & , m = k - 1 \\ 0 & , \text{else} \end{cases}$$



a transition probability on S is defined, that is proportional to the number of molecules in container A . For this transition matrix the binomial distribution on S ,

$$\mu(k) := \binom{r}{k} 2^{-r} \quad (k \in S \equiv \{0, 1, \dots, r\}) ,$$

is an invariant measure.

PROOF. Since μ is a probability measure, we only have to show that $\mu(k) = \sum_{m=0}^r p(m, k) \mu(m)$ for $k = 0, 1, \dots, r$ is valid. With $k = 1, \dots, r - 1$ we have

$$\begin{aligned} \sum_{m=0}^r p(m, k) \mu(m) &= p(k+1, k) \mu(k+1) + p(k-1, k) \mu(k-1) \\ &\equiv 2^{-r} \left[\binom{r}{k+1} \frac{k+1}{r} + \binom{r}{k-1} \frac{r-(k-1)}{r} \right] \\ &= 2^{-r} \left[\frac{(r-1)!}{k! (r-(k+1))!} + \frac{(r-1)!}{(k-1)! (r-k)!} \right] \\ &= 2^{-r} \frac{(r-1)!}{(k-1)! (r-k-1)!} \left[\frac{1}{k} + \frac{1}{r-k} \right] \\ &= 2^{-r} \frac{r!}{k! (r-k)!} \\ &\equiv \mu(k) . \end{aligned}$$

In the cases $k = 0$ and $k = r$ only one summand does not vanish. □

Now we show how to associate with each class of recurrent states a stationary measure; the Markov chain decouples on these classes. We constantly use

$$P_x(X_n = y) = p^n(x, y) \quad (x, y \in S; n \in \mathbb{N}) ,$$

where $p^n(x, y)$ is the n -fold matrix product.

Theorem 2.3. *Let x be recurrent and $T \equiv T_x := \inf\{n \in \mathbb{N} : X_n = x\}$ its first hitting time. Then*

$$\mu(y) := \mathbb{E}_x \left(\sum_{n=0}^{T-1} 1_{\{X_n=y\}} \right) = \sum_{n=0}^{\infty} P_x(X_n = y, T > n) \quad (y \in S)$$

defines a stationary measure.

PROOF. First we prove the equation $\mu p = \mu$; this way we prove that $\mu(y) < \infty$ for all $y \in S$. Note that $\mu(x) = 1$.

(a) $\sum_{y \in S} \mu(y) p(y, z) = \mu(z)$ for all $z \in S$:

1) If $z \neq x$, the Markov property (Thm 1.7) implies:

$$\begin{aligned} \sum_{y \in S} \mu(y) p(y, z) &\stackrel{\text{Fubini}}{=} \sum_{n=0}^{\infty} \sum_{y \in S} P_x(X_n = y, T > n) \cdot P_y(X_1 = z) \\ &\stackrel{\text{MP}}{=} \sum_{n=0}^{\infty} \sum_{y \in S} P_x(X_n = y, T > n, X_{n+1} = z) \\ &= \sum_{n=0}^{\infty} P_x(T > n, X_{n+1} = z) \\ &\stackrel{z \neq x}{=} \sum_{n=0}^{\infty} P_x(T > n+1, X_{n+1} = z) \\ &= \sum_{n=1}^{\infty} P_x(T > n, X_n = z) \\ &\stackrel{z \neq x}{=} \sum_{n=0}^{\infty} P_x(T > n, X_n = z) \quad \equiv \quad \mu(z). \end{aligned}$$

2) If $z = x$, again by the Markov property 1.7:

$$\begin{aligned} \sum_{y \in S} \mu(y) p(y, x) &\stackrel{\text{ME}}{=} \sum_{n=0}^{\infty} \sum_{y \in S} P_x(X_n = y, T > n, X_{n+1} = x) \\ &= \sum_{n=0}^{\infty} P_x(T = n+1) = \rho_{xx} \stackrel{x \text{ rec.}}{=} 1 = \mu(x). \end{aligned}$$

(b) $\mu(y) < \infty$ for all $y \in S$:

1) In case $\rho_{xy} > 0$: By iteration of (a) we get: $\mu = \mu p^n$ for $n \in \mathbb{N}$ so

$$1 = \mu(x) \stackrel{(a)}{=} (\mu p^n)(x) = \sum_{y \in S} \mu(y) p^n(y, x) \quad (n \in \mathbb{N}).$$

Consequently necessarily $\mu(y) < \infty$, if $p^n(y, x) > 0$ with some $n \in \mathbb{N}$; since $p^n(y, x) = P_y(X_n = x)$, the latter is implied by $\rho_{yx} \equiv P_y(T_x < \infty) > 0$, which in the case considered $\rho_{xy} > 0$ by recurrence of x follows from Thm 1.11 (hence $x \sim y$).

2) If $\rho_{xy} = 0$, the definition of μ gives $\mu(y) = 0 (< \infty)$. □

Theorem 2.4 (Uniqueness of stationary measures). *Let $(X_n)_{n \in \mathbb{N}_0}$ be irreducible, i.e. S has only one class of recurrent states. Then the stationary measure μ from Theorem 2.3 is unique up to multiplication by constants.*

PROOF. Let $a \in S$ be a recurrent state and μ the stationary measure belonging to a according to 2.3. If ν denotes a further stationary measure, we have to show:

$$\nu(z) = \mu(z) \cdot \nu(a) \quad (z \in S).$$

By stationarity of ν we obtain iteratively for $z \in S$:

$$\begin{aligned} \nu(z) &= \sum_{y \in S} \nu(y) p(y, z) \\ &= \nu(a) p(a, z) + \sum_{y \neq a} \nu(y) p(y, z) \\ &= \nu(a) p(a, z) + \sum_{y \neq a} \left(\sum_{x \in S} \nu(x) p(x, y) \right) p(y, z) \\ &= \nu(a) p(a, z) + \sum_{y \neq a} \nu(a) p(a, y) p(y, z) + \sum_{y \neq a} \sum_{x \neq a} \nu(x) p(x, y) p(y, z) \\ &= \nu(a) P_a(X_1 = z) + \sum_{y \neq a} \nu(a) P_a(X_1 \neq a, X_2 = z) \\ &\quad + P_\nu(X_0 \neq a, X_1 \neq a, X_2 = z) \\ &= \dots = \\ &= \nu(a) \sum_{m=1}^n P_a(X_k \neq a \text{ for } 1 \leq k < m, X_m = z) \\ &\quad + P_\nu(X_0 \neq a, X_1 \neq a, \dots, X_{n-1} \neq a, X_n = z) \\ &\geq \nu(a) \cdot \mu(z) \end{aligned}$$

($n \rightarrow \infty$) by definition of μ ; therefore for $n \in \mathbb{N}$:

$$\nu(a) = \sum_{z \in S} \nu(z) p^n(z, a) \geq \nu(a) \sum_{z \in S} \mu(z) p^n(z, a) = \nu(a) \mu(a) = \nu(a).$$

In the previous estimate $\nu(z) \geq \nu(a) \mu(z)$ the inequality " $>$ " can only be valid if $p^n(z, a) = 0$ for each $n \in \mathbb{N}$. By irreducibility of for any z there exists $n \in \mathbb{N}$ with $p^n(z, a) > 0$. Therefore $\nu(z) = \nu(a) \mu(z)$. \square

We give a necessary condition for the normability of stationary measures:

Satz 2.5. *If there exists an invariant measure μ , all states y with $\mu(y) > 0$ are recurrent.*

PROOF. For $n \in \mathbb{N}$ we have by stationarity $\mu = \mu p^n$, hence by Fubini

$$\sum_{n=1}^{\infty} \mu(y) = \sum_{x \in S} \mu(x) \sum_{n=1}^{\infty} p^n(x, y) \stackrel{1.10}{=} \sum_{x \in S} \mu(x) \frac{\rho_{xy}}{1 - \rho_{yy}} \leq \frac{\mu(S)}{1 - \rho_{yy}}.$$

By hypothesis $\sum_{n=1}^{\infty} \mu(y) = \infty$ and $\mu(S) = 1 < \infty$, hence $\rho_{yy} = 1$. \square

Theorem 2.6. *Let $(X_n)_{n \in \mathbb{N}_0}$ be irreducible and μ an invariant measure. Then*

$$\mu(x) = \frac{1}{\mathbb{E}_x(T_x)} \quad (x \in S).$$

PROOF. Note first that all elements of S are recurrent: Each element with positive mass w.r.t. μ is recurrent by 2.5; but since X is irreducible, this recurrence transfers to all other elements.

Consequently for each fixed $x \in S$ by 2.3 there exists a stationary measure μ_0 :

$$\mu_0(z) \equiv \sum_{n \in \mathbb{N}_0} P_x(X_n = z, T_x > n) \quad \text{and} \quad \mu_0(x) = 1.$$

Consequently by Fubini:

$$\sum_{z \in S} \mu_0(z) = \sum_{n=0}^{\infty} \sum_{z \in S} P_x(X_n = z, T_x > n) = \sum_{n=0}^{\infty} P_x(T_x > n) = \mathbb{E}_x(T_x).$$

By the uniqueness statement in 2.4 this means for the normed measure μ :

$$\mu(y) = \frac{\mu_0(y)}{\sum_{z \in S} \mu_0(z)} = \frac{\mu_0(y)}{\mathbb{E}_x(T_x)} \quad (y \in S),$$

whence by $y = x$ the claim follows, since $\mu_0(x) = 1$. \square

$x \in S$ is called *positively recurrent*, if $\mathbb{E}_x(T_x) < \infty$; in the other case x is called *null recurrent*.

”Positively recurrent” is stronger than ”recurrent”. Positive and null recurrence are properties of classes. In the Ehrenfest model 2.2 every state is positively recurrent.

Corollary 2.7. *Let $(X_n)_{n \in \mathbb{N}_0}$ be irreducible. Then the following statements are equivalent:*

- i) *There exists an invariant measure;*
- ii) *There exists a positively recurrent state;*
- iii) *All states are positively recurrent.*

PROOF. iii) \Rightarrow ii) trivial.

ii) \Rightarrow i) Let x be positively recurrent. By 2.3 there exists a stationary measure μ_0 with total mass $\mu_0(S) = \sum_{z \in S} \mu_0(z) = \mathbb{E}_x(T_x)$ (proof of 2.6), which by positive recurrence is finite. The norming factor μ is therefore invariant:

$$\mu(y) := \frac{\mu_0(y)}{\mathbb{E}_x(T_x)} \equiv \frac{1}{\mathbb{E}_x(T_x)} \sum_{n \in \mathbb{N}_0} P_x(X_n = y, T_x > n) \quad (y \in S).$$

i) \Rightarrow iii) Let μ be the invariant measure. By irreducibility $\mu(x) > 0$ for all $x \in S$ (every state x is recurrent, so that $\mu_0(x) = 1$ for the stationary measure μ_0 given according to 2.3; by 2.4 we therefore must have $\mu(x) > 0$). From 2.6 we conclude: $\mathbb{E}_x(T_x) = \frac{1}{\mu(x)} < \infty$ for each $x \in S$. \square

We now discuss criteria under which p^n converges to the invariant measure.

Example 2.8. On $S := \{1, 2\}$ $p := \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix}$ defines a transition matrix. We have

$$p^{2n} = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix} \quad \text{and} \quad p^{2n+1} = \begin{pmatrix} 0 & 1 \\ 1 & 0 \end{pmatrix} \equiv p \quad (n \in \mathbb{N}).$$

In this case no convergence of $p^n(x, y)$ is given.

Periodicity prevents convergence to the invariant measure.

Definition 2.9. For a recurrent $x \in S$ let²

$$I_x := \{n \in \mathbb{N}_0 : p^n(x, x) > 0\}.$$

$d_x := \text{gcd}(I_x)$ is called *period* of x .

By the Chapman-Kolmogorov equation I_x is a semigroup.

In the above example 2.8 we have $I_1 = I_2 = \{\text{gerade Zahlen}\}$ and $d_1 = d_2 = 2$.

Lemma 2.10. Let $x, y \in S$ be recurrent with $x \sim y$. Then $d_x = d_y$.

PROOF. We show³: $d_y \mid d_x$. Since our arguments are symmetric in x and y , this implies the claim, for by switching the roles of x and y we also have $d_x \mid d_y$. Wlog we may assume $x \neq y$. By the equivalence $x \sim y$ we therefore have $\rho_{xy} > 0$ and $\rho_{yx} > 0$; in particular there exist $m, n \in \mathbb{N}$ with $p^m(x, y) > 0$ and $p^n(y, x) > 0$. By Chapman-Kolmogorov this implies

$$p^{n+m}(y, y) \geq p^n(y, x) p^m(x, y) > 0.$$

Hence by the above definition we obtain $d_y \mid n + m$.

Let now an arbitrary $k \in I_x$ be given. By what has just been proved $d_y \mid n + m$ we only have to show that also $d_y \mid n + m + k$. These two statements imply $d_y \mid k$ and thus the claim. By Chapman-Kolmogorov and $k \in I_x$ we get

$$p^{n+k+m}(y, y) \geq p^n(y, x) p^k(x, x) p^m(x, y) > 0,$$

and thus $d_y \mid n + k + m$. □

Definition 2.11. (a) A state $x \in S$ is called *aperiodic*, if $d_x = 1$ holds.

(b) An irreducible, recurrent Markov chain is called *aperiodic*, if each state is aperiodic.

As indicated in the above example, we shall see that aperiodicity is a criterion for the convergence of the transition probabilities to the invariant measure. The proof of this fact is prepared by the following lemma.

Lemma 2.12. For aperiodic x there exists $m_0 \in \mathbb{N}$ with $p^m(x, x) > 0$ for all $m \geq m_0$.

²Reminder: $p^n(x, y) \equiv P_x(X_n = y)$ for $x, y \in S$ and $n \in \mathbb{N}_0$.

³As usual "|" abbreviates "is a divisor of".

PROOF. We first prove that there is $N \in \mathbb{N}$ such that $N, N + 1 \in I_x$. For this purpose let $n_0, n_0 + k \in I_x$ be fixed. In case $k = 1$ the proof is finished. In case $k \geq 2$ we choose $n_1 \in I_x$ with $k \nmid n_1$ (since $d_x = 1$). For this we have (division with remainder)

$$n_1 = mk + r_1 \quad (m \in \mathbb{N}_0, 0 < r_1 < k)$$

and by the semigroup property of I_x

$$(m + 1)(n_0 + k) \in I_x \quad \text{and} \quad (m + 1)n_0 + n_1 \in I_x.$$

For these two elements we have:

$$\begin{aligned} \left| (m + 1)(n_0 + k) - ((m + 1)n_0 + n_1) \right| &= |(m + 1)k - n_1| \\ &\equiv |(m + 1)k - (mk + r_1)| = k - r_1 < k. \end{aligned}$$

If $k - r_1 = 1$, the claim holds with $N := (m + 1)n_0 + n_1$. If $k - r_1 > 1$, we repeat the step performed with $\tilde{n}_0 := (m + 1)n_0 + n_1$ and $\tilde{k} := k - r_1$. After finitely many iterations we obtain $N \in \mathbb{N}$ with $N, N + 1 \in I_x$.

With this the claim of the Lemma follows with $m_0 := N^2$, since for $m \geq m_0$ we have

$$m - N^2 = kN + r \quad (k \in \mathbb{N}_0, 0 \leq r < N)$$

(division with remainder), so that

$$m = N^2 + kN + r = (N - r + k)N + r(1 + N) \in I_x$$

by the semigroup property of I_x . □

Theorem 2.13 (Invariant measure is limit of transition probabilities). *Let the Markov chain $(X_n)_{n \in \mathbb{N}_0}$ be aperiodic and possess the invariant measure μ . Then*

$$p^n(x, y) \xrightarrow{n \rightarrow \infty} \mu(y) = \frac{1}{\mathbb{E}_y(T_y)} \quad (x, y \in S).$$

PROOF(COUPLING OF PROCESSES, W. DÖBLIN). On $S^2 \equiv S \times S$ setting

$$q((x_1, y_1), (x_2, y_2)) := p(x_1, x_2)p(y_1, y_2) \quad (x_1, x_2, y_1, y_2 \in S)$$

defines a transition probability. Let $(X_n, Y_n)_{n \in \mathbb{N}_0}$ be the canonical Markov chain associated with q , that is the Markov chain with state space S^2 on

$$(\Omega, \mathcal{F}, \mathbb{P}) := \left((S^2)^{\mathbb{N}_0}, (\mathcal{S}^2)^{\mathbb{N}_0}, P_\varrho \right),$$

where P_ϱ is the probability measure related to q and an initial distribution ϱ (on $\mathcal{S}^2 \equiv \mathcal{S} \otimes \mathcal{S}$) according to Kolmogorov.

With 2.12 we now prove irreducibility of $(X_n, Y_n)_{n \in \mathbb{N}_0}$; from this we get that this coupled process hits the diagonal of S^2 in finite time. This will imply convergence.

1) $(X_n, Y_n)_{n \in \mathbb{N}_0}$ is irreducible: If $x_1, x_2, y_1, y_2 \in S$ are fixed, irreducibility of X provides times $k, l \in \mathbb{N}$ with

$$p^k(x_1, x_2) > 0 \quad \text{and} \quad p^l(y_1, y_2) > 0.$$

Aperiodicity yields according to 2.12 also an $m_0 \in \mathbb{N}$, such that for $m \geq m_0$ we have

$$p^{m+l}(x_2, x_2) > 0 \quad \text{and} \quad p^{m+k}(y_2, y_2) > 0.$$

Hence by Chapman-Kolmogorov also

$$\begin{aligned} q^{k+l+m}((x_1, y_1), (x_2, y_2)) \\ &\equiv p^{k+l+m}(x_1, x_2) p^{k+l+m}(y_1, y_2) \\ &\geq p^k(x_1, x_2) p^{m+l}(x_2, x_2) p^l(y_1, y_2) p^{m+k}(y_2, y_2) > 0. \end{aligned}$$

Therefore S^2 consists of a unique equivalence class. For irreducibility we have to show that all states in S^2 are recurrent. By 2.5 for this we need a q -invariant measure ν with $\nu(x, y) > 0$ for all $(x, y) \in S^2$. But setting

$$\nu(x, y) := \mu(x) \mu(y) \quad (x, y \in S)$$

provides a q -invariant measure on S^2 by p -invariance of μ :

$$\begin{aligned} \sum_{(x_1, x_2) \in S^2} \nu(x_1, x_2) q((x_1, x_2), (y_1, y_2)) &\equiv \sum_{(x_1, x_2)} \mu(x_1) \mu(x_2) p(x_1, y_1) p(x_2, y_2) \\ &= \sum_{x_1} \mu(x_1) p(x_1, y_1) \sum_{x_2} \mu(x_2) p(x_2, y_2) = \mu(y_1) \mu(y_2) \equiv \nu(y_1, y_2) \end{aligned}$$

for $(y_1, y_2) \in S^2$; moreover $\nu(y_1, y_2) \equiv \mu(y_1) \mu(y_2) \stackrel{2.6}{=} \frac{1}{\mathbb{E}_{y_1}(T_{y_1})} \frac{1}{\mathbb{E}_{y_2}(T_{y_2})} \stackrel{2.7 \text{ iii)}}{>} 0$.

2) Denote by T the first hitting time of the diagonal $D := \{(x, x) : x \in S\}$,

$$T := \inf\{n \in \mathbb{N} : (X_n, Y_n) \in D\},$$

$T_{(x,x)}$ the time of first visit in $(x, x) \in D$. Then on the one hand $T \leq T_{(x,x)}$. If ϱ is an arbitrary initial distribution on S^2 , on the other hand by the recurrence proved in 1) we get $T_{(x,x)} < \infty$ P_ϱ -a.s.; in particular $T < \infty$ P_ϱ -a.s. .

X_n and Y_n possess on $\{T \leq n\}$ identical laws ($n \in \mathbb{N}$), since for $y \in S$:

$$\begin{aligned} P_\varrho(X_n = y, T \leq n) &= \sum_{m=1}^n P_\varrho(T = m, X_n = y) \\ &= \sum_{m=1}^n \sum_{x \in S} P_\varrho(T = m, X_m = x, X_n = y) = \end{aligned}$$

$$\begin{aligned}
&= \sum_{m=1}^n \sum_{x \in S} P_\varrho(X_n = y \mid T = m, X_m = x) P_\varrho(T = m, X_m = x) \\
&\stackrel{\text{ME}}{=} \sum_{m=1}^n \sum_{x \in S} P_\varrho(X_n = y \mid X_m = x) P_\varrho(T = m, X_m = x) \\
&= \sum_{m=1}^n \sum_{x \in S} P_\varrho(Y_n = y \mid Y_m = x) P_\varrho(T = m, Y_m = x) \\
&= \dots \dots \dots \stackrel{\text{same arg.}}{=} P_\varrho(Y_n = y, T \leq n) .
\end{aligned}$$

Here we used that X and Y possess identical transition probability p .

3) Now we prove the claim of the theorem; for this purpose we show the following (stronger) convergence:

$$\sum_{y \in S} |p^n(x, y) - \mu(y)| \xrightarrow{n \rightarrow \infty} 0$$

for all $x \in S$; the equality $\mu(y) = 1/\mathbb{E}_y(T_y)$ is already clear by 2.6.

For an arbitrary $x \in S$ we fix the initial distribution

$$\varrho := \delta_x \otimes \mu$$

on S^2 for the coupled process. Thus for all $y \in S$

$$\begin{aligned}
p^n(x, y) &= P_\varrho(X_n = y) \\
&= P_\varrho(X_n = y, T \leq n) + P_\varrho(X_n = y, T > n) \\
&\stackrel{2)}{=} P_\varrho(Y_n = y, T \leq n) + P_\varrho(X_n = y, T > n)
\end{aligned}$$

by equality of the laws proven in 2), and

$$\mu(y) = P_\varrho(Y_n = y) \equiv P_\varrho(Y_n = y, T \leq n) + P_\varrho(Y_n = y, T > n)$$

by the p -invariance of μ ; in summary

$$\begin{aligned}
\sum_{y \in S} |p^n(x, y) - \mu(y)| &= \sum_{y \in S} |P_\varrho(X_n = y) - P_\varrho(Y_n = y)| \\
&= \sum_{y \in S} |P_\varrho(X_n = y, T > n) - P_\varrho(Y_n = y, T > n)| \\
&\leq \sum_{y \in S} [P_\varrho(X_n = y, T > n) + P_\varrho(Y_n = y, T > n)] \\
&= 2 P_\varrho(T > n) \xrightarrow{n \rightarrow \infty} 0,
\end{aligned}$$

since T is P_ϱ -a.s. finite, as seen in 2). □

3. Stationary Processes

In this chapter we consider stochastic processes $X = (X_n)_{n \in \mathbb{N}_0}$ on a fixed probability space $(\Omega, \mathcal{F}, \mathbb{P})$ with values in a Polish space S (equipped with the Borel σ -algebra $\mathcal{S} := \mathcal{B}(S)$). This family of \mathcal{F} - \mathcal{S} -measurable maps can also be considered as a random sequence

$$X : \Omega \longrightarrow S^{\mathbb{N}_0}, \quad \omega \mapsto (X_n(\omega))_{n \in \mathbb{N}_0},$$

which is \mathcal{F} - $\mathcal{S}^{\mathbb{N}_0}$ -measurable, with product σ -algebra

$$\mathcal{S}^{\mathbb{N}_0} := \sigma \left(\bigcup_{n \in \mathbb{N}_0} \pi_{\{n\}}^{-1}[B_n] : B_n \in \mathcal{S} \right) = \sigma \left(\bigcup_{n \in \mathbb{N}_0} \pi_n^{-1}[B] : B \in \mathcal{S}^{n+1} \right);$$

where the second generating system is \cap -stable, in contrast to the first. The measure defined by

$$P_X \equiv P_{(X_n)_{n \in \mathbb{N}_0}} := \mathbb{P} \circ X^{-1}$$

on $\mathcal{S}^{\mathbb{N}_0}$ is the *law of X*.

If only distribution properties are relevant, instead of X we can wlog also study its *canonical representation* $(Y)_n := (\pi_{\{n\}})_{n \in \mathbb{N}_0}$ on $(S^{\mathbb{N}_0}, \mathcal{S}^{\mathbb{N}_0}, P_X)$.

Definition 3.1. A stochastic process $X = (X_n)_{n \in \mathbb{N}_0}$ is called *stationary*, if we have:

$$P_{(X_n)_{n \in \mathbb{N}_0}} = P_{(X_{n+k})_{n \in \mathbb{N}_0}} \quad (\forall k \in \mathbb{N}).$$

The distribution of a stationary process does not "move"; this will be enforced in the following Lemma:

Lemma 3.2. $X = (X_n)_{n \in \mathbb{N}_0}$ is stationary iff we have:

$$P_{(X_0, \dots, X_n)} = P_{(X_k, \dots, X_{k+n})} \quad (k \in \mathbb{N}, n \in \mathbb{N}_0).$$

PROOF. " \Rightarrow " For all $k \in \mathbb{N}$, $n \in \mathbb{N}_0$ and $B \in \mathcal{S}^{n+1}$ we have:

$$\begin{aligned} P_{(X_0, \dots, X_n)}(B) &\equiv \mathbb{P}\{(X_0, \dots, X_n) \in B\} \\ &= \mathbb{P}\{(X_m)_{m \in \mathbb{N}_0} \in \pi_n^{-1}(B)\} \\ &\stackrel{\text{stat}}{=} \mathbb{P}\{(X_{m+k})_{m \in \mathbb{N}_0} \in \pi_n^{-1}(B)\} \\ &= \mathbb{P}\{(X_k, \dots, X_{k+n}) \in B\} \equiv P_{(X_k, \dots, X_{k+n})}(B). \end{aligned}$$

" \Leftarrow " By hypothesis we have for all $k \in \mathbb{N}$, $n \in \mathbb{N}_0$ and $B \in \mathcal{S}^{n+1}$:

$$P_{(X_m)_{m \in \mathbb{N}_0}}(\pi_n^{-1}(B)) = P_{(X_{m+k})_{m \in \mathbb{N}_0}}(\pi_n^{-1}(B))$$

(see calculation above). But since $\{\bigcup_{n \in \mathbb{N}_0} \pi_n^{-1}(B) : B \in \mathcal{S}^{n+1}\}$ is a \cap -stable generator of $\mathcal{S}^{\mathbb{N}_0}$, this implies $P_{(X_m)_{m \in \mathbb{N}_0}} = P_{(X_{m+k})_{m \in \mathbb{N}_0}}$ by the uniqueness theorem for measures. \square

Example 3.3 (Markov chain with transition probability p). Let $(X_n)_{n \in \mathbb{N}_0}$ be a Markov chain on a countable space S (equipped with $\mathcal{S} := \mathcal{B}(S) \equiv \mathfrak{P}(S)$) with transition probability p and invariant measure μ . Then $(X_n)_{n \in \mathbb{N}_0}$ is stationary on $(\Omega, \mathcal{F}, \mathbb{P}) := (S^{\mathbb{N}_0}, \mathcal{S}^{\mathbb{N}_0}, P_\mu)$.

PROOF. Note first that we have for all $B := B_0 \times B_1 \times \cdots \times B_n \in \mathcal{S}^{n+1}$:

$$\begin{aligned}
(P_\mu)_{(X_1, \dots, X_{n+1})}(B) &\equiv P_\mu\{X_1 \in B_0, X_2 \in B_1, \dots, X_{n+1} \in B_n\} \\
&= P_\mu\{X_0 \in S, X_1 \in B_0, X_2 \in B_1, \dots, X_{n+1} \in B_n\} \\
&= \sum_{z \in S} \mu(z) \sum_{x_0 \in B_0} p(z, x_0) \sum_{x_1 \in B_1} p(x_0, x_1) \cdots \sum_{x_n \in B_n} p(x_{n-1}, x_n) \\
&= \sum_{x_0 \in B_0} \sum_{z \in S} \mu(z) p(z, x_0) \sum_{x_1 \in B_1} p(x_0, x_1) \cdots \sum_{x_n \in B_n} p(x_{n-1}, x_n) \\
&\stackrel{\text{inv}}{=} \sum_{x_0 \in B_0} \mu(x_0) \sum_{x_1 \in B_1} p(x_0, x_1) \cdots \sum_{x_n \in B_n} p(x_{n-1}, x_n) \\
&= P_\mu\{X_0 \in B_0, X_1 \in B_1, \dots, X_n \in B_n\} \\
&= (P_\mu)_{(X_0, X_1, \dots, X_n)}(B).
\end{aligned}$$

By k -fold iteration of this argument we obtain the criterion for stationarity from 3.2. \square

Example 3.4 (Rotation on circle). Let $(\Omega, \mathcal{F}, \mathbb{P}) := ([0, 1), \mathcal{B}[0, 1), \lambda|_{\mathcal{F}})$, where λ denotes the Lebesgue measure. Then for each fixed $\theta \in [0, 1)$ the process $(X_n)_{n \in \mathbb{N}_0}$,

$$X_n : \Omega \longrightarrow S := \Omega, \quad X_n(\omega) := \omega + n \cdot \theta \pmod{1}, \quad n \in \mathbb{N}_0,$$

is a stationary Markov chain on $(S^{\mathbb{N}_0}, \mathcal{S}^{\mathbb{N}_0}, P_\lambda)$ with respect to the transition probability

$$p : S \times \mathcal{S} \longrightarrow [0, 1], \quad p(x, A) := \begin{cases} 1, & \text{if } y = x + \theta \pmod{1} \in A, \\ 0, & \text{else.} \end{cases}$$

PROOF. By translation invariance of the Lebesgue measure λ is p -invariant, since for $A \in \mathcal{S}$

$$\int_0^1 \lambda(dz) p(z, A) = \lambda(A - \theta \pmod{1}) = \lambda(A).$$

Hence as in Example 3.3 for all $B := B_0 \times B_1 \times \cdots \times B_n \in \mathcal{S}^{n+1}$:

$$\begin{aligned}
(P_\lambda)_{(X_1, \dots, X_{n+1})}(B) &= P_\lambda\{X_0 \in S, X_1 \in B_0, X_2 \in B_1, \dots, X_{n+1} \in B_n\} \\
&= \int_\Omega \lambda(dz) \int_{B_0} p(z, dx_0) \int_{B_1} p(x_0, dx_1) \cdots \int_{B_n} p(x_{n-1}, dx_n) \\
&= \int_{B_0} \int_\Omega \lambda(dz) p(z, dx_0) \int_{B_1} p(x_0, dx_1) \cdots \int_{B_n} p(x_{n-1}, dx_n) \\
&\stackrel{\text{inv}}{=} \int_{B_0} \lambda(dx_0) \int_{B_1} p(x_0, dx_1) \cdots \int_{B_n} p(x_{n-1}, dx_n) \\
&= (P_\lambda)_{(X_0, X_1, \dots, X_n)}(B),
\end{aligned}$$

and thus stationarity again from Lemma 3.2. \square

Theorem 3.5. *Let the process $(X_n)_{n \in \mathbb{N}_0}$ with Polish state space (S, \mathcal{S}) be stationary and let $g : S^{\mathbb{N}_0} \rightarrow S'$ be $\mathcal{S}^{\mathbb{N}_0}$ - \mathcal{S}' -measurable, where (S', \mathcal{S}') is also Polish. Then*

$$Y_k := g(X_k, X_{k+1}, \dots) \quad (k \in \mathbb{N}_0)$$

is stationary (in S').

PROOF. By measurability of g for each $k \in \mathbb{N}_0$ also

$$g_k : S^{\mathbb{N}_0} \rightarrow S', \quad x \mapsto g \circ \theta_k(x),$$

is measurable, where $\theta \equiv (\theta_k)_{k \in \mathbb{N}_0}$ (see 1.5) denotes the measurable shift

$$\theta_k : S^{\mathbb{N}_0} \rightarrow S^{\mathbb{N}_0}, \quad (x_n)_n \mapsto (x_{n+k})_n.$$

Let now $B \in (\mathcal{S}')^{\mathbb{N}_0}$ be fixed; by measurability of all g_k also $A := (g_0, g_1, \dots)^{-1}(B)$ is measurable and by $Y_k = g_k((X_n)_n)$ we obtain for $m \in \mathbb{N}$:

$$\begin{aligned} \mathbb{P}_{(Y_k)_{k \in \mathbb{N}_0}}(B) &\equiv \mathbb{P}((Y_k)_{k \in \mathbb{N}_0} \in B) &= \mathbb{P}((X_n)_{n \in \mathbb{N}_0} \in A) \\ &\stackrel{X \text{ stat}}{\equiv} \mathbb{P}((X_{n+m})_{n \in \mathbb{N}_0} \in A) &= \mathbb{P}((Y_{k+m})_{k \in \mathbb{N}_0} \in B) \\ &\equiv \mathbb{P}_{(Y_{k+m})_{k \in \mathbb{N}_0}}(B), \end{aligned}$$

hence the stationarity of Y . □

Example 3.6 (Bernoulli-Shift). On $(\Omega, \mathcal{F}, \mathbb{P}) := ([0, 1), \mathcal{B}[0, 1), \lambda|_{\mathcal{F}})$ $(Y_n)_{n \in \mathbb{N}_0}$,

$$Y_n : \Omega \rightarrow \Omega, \quad Y_n := \begin{cases} \text{id}_\Omega, & n = 0, \\ 2Y_{n-1} \pmod{1}, & n \in \mathbb{N}, \end{cases}$$

is stationary.

PROOF. Let $(X_n)_{n \in \mathbb{N}_0}$ be a Bernoulli sequence with rate $\frac{1}{2}$, realized as product measure $\tilde{\mathbb{P}}$ on $\tilde{\Omega} := \{0, 1\}^{\mathbb{N}_0}$; hence $(X_n)_n$ is a sequence of i.i.d. random variables in $S := \{0, 1\}$ with $\tilde{\mathbb{P}}\{X_n = 0\} = \tilde{\mathbb{P}}\{X_n = 1\} = \frac{1}{2}$. Then $(X_n)_n$ is stationary. Moreover

$$g : \tilde{\Omega} \equiv \{0, 1\}^{\mathbb{N}_0} \longrightarrow \Omega \equiv [0, 1), \quad (x_n)_n \mapsto \sum_{n=0}^{\infty} x_n 2^{-n-1} \pmod{1}$$

is measurable and so $\tilde{\mathbb{P}} \circ g^{-1} = \mathbb{P}$ (dyadic intervals may be written as sets of the form $\{X_0 = i_0, \dots, X_k = i_k\}$ with $i_0, \dots, i_k \in \{0, 1\}$). Because of Theorem 3.5 we have that

$$Z_k := g(X_k, X_{k+1}, \dots) \quad (k \in \mathbb{N}_0)$$

is stationary; on the other hand we have:

$$\begin{aligned} 2Z_0 \equiv 2g(X_0, X_1, \dots) &= X_0 + \sum_{n=1}^{\infty} X_n 2^{-n} \pmod{1} \\ &= X_0 + \sum_{n=0}^{\infty} X_{n+1} 2^{-(n+1)} \pmod{1} \\ &= g(X_1, X_2, \dots) \equiv Z_1; \end{aligned}$$

by iteration we obtain: $2Z_{n-1} = Z_n$ ($n \in \mathbb{N}$), hence with Z also Y is stationary. \square

Definition 3.7 (measure preserving map). Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space. A \mathcal{F} - \mathcal{F} -measurable mapping $\varphi : \Omega \rightarrow \Omega$ is called *measure preserving*, if we have: $\mathbb{P} \circ \varphi^{-1} = \mathbb{P}$.

Remark 3.8. Let φ be measure preserving on $(\Omega, \mathcal{F}, \mathbb{P})$ and $X : \Omega \rightarrow S$ a \mathcal{F} - \mathcal{S} -measurable map with values in a Polish space (S, \mathcal{S}) . Then $(X_n)_{n \in \mathbb{N}_0}$ with

$$X_n := \begin{cases} X, & n = 0 \\ X \circ \varphi^n, & n \in \mathbb{N} \end{cases}$$

is stationary.

PROOF. For $B \in \mathcal{S}^{n+1}$ we have:

$$\begin{aligned} \mathbb{P}_{(X_0, \dots, X_n)}(B) &\equiv \mathbb{P}((X_0, \dots, X_n) \in B) \\ &\stackrel{\varphi \text{ m.p.}}{=} \mathbb{P}((X_0, \dots, X_n) \circ \varphi^k \in B) = \mathbb{P}_{(X_k, \dots, X_{k+n})}(B), \end{aligned}$$

so that stationarity follows from Lemma 3.2. \square

The situation of the preceding remark does not only provide an example for a stationary sequence. It already depicts the general situation.

Satz 3.9 (standard model for stationary sequences).

Let $(Y_n)_{n \in \mathbb{N}_0}$ be stationary on $(\Omega, \mathcal{F}, \mathbb{P})$ with values in a Polish space (S, \mathcal{S}) . Then there exists a probability space $(\Omega', \mathcal{F}', \mathbb{P}')$ with a measure preserving map $\varphi : \Omega' \rightarrow \Omega'$ and a random variable $X_0 : \Omega' \rightarrow S$ such that with $X_n := X_0 \circ \varphi^n$ ($n \in \mathbb{N}$) we have:

$$\mathbb{P}'_{(X_n)_{n \in \mathbb{N}_0}} = \mathbb{P}_{(Y_n)_{n \in \mathbb{N}_0}}.$$

PROOF. Let $(\Omega', \mathcal{F}', \mathbb{P}') := (S^{\mathbb{N}_0}, \mathcal{S}^{\mathbb{N}_0}, P_{(Y_n)_{n \in \mathbb{N}_0}})$ and $X_0 := \pi_{\{0\}}$ (projection at time 0) and $\varphi := \theta_1$ (shift). By stationarity of Y φ is measure preserving, since for $A' \in \mathcal{F}'$ we have:

$$\begin{aligned} \mathbb{P}'(\varphi^{-1}(A')) &= \mathbb{P}((Y_n)_n \in \varphi^{-1}(A')) \\ &= \mathbb{P}((Y_{n+1})_n \in A') \\ &\stackrel{Y \text{ stat}}{=} \mathbb{P}((Y_n)_n \in A') = \mathbb{P}'(A'). \end{aligned}$$

The claimed equality of laws follows from the definition of \mathbb{P}' . \square

Definition 3.10 (invariant, ergodic). Let φ be a measure preserving mapping on $(\Omega, \mathcal{F}, \mathbb{P})$.

$A \in \mathcal{F}$ is called $\left\{ \begin{array}{l} \text{invariant} \\ \text{strictly invariant} \end{array} \right\}$, if $\left\{ \begin{array}{l} \varphi^{-1}(A) = A \text{ } \mathbb{P}\text{-a.s.} \\ \varphi^{-1}(A) = A \end{array} \right\}$.

φ is called *ergodic*, if for all $A \in \mathcal{I} := \{ \text{invariant sets} \}$ we have: $\mathbb{P}(A) \in \{0, 1\}$.

Remark 3.11. i) \mathcal{I} is a σ -algebra (sub- σ -algebra of \mathcal{F});

ii) For $A \in \mathcal{I}$ there exists a strictly invariant set $B \in \mathcal{F}$ with $B = A$ \mathbb{P} -a.s.
(for example $B := \liminf_{n \rightarrow \infty} \varphi^{-n}(A)$).

iii) For $A \in \mathcal{I}$ there exists $B \in \mathcal{I} := \bigcap_{n=1}^{\infty} \sigma(X_n, X_{n+1}, \dots)$ with⁴ $B = A$ \mathbb{P} -a.s.
(for example again $B := \liminf_{n \rightarrow \infty} \varphi^{-n}(A)$, since $B = \varphi^{-k}(B) \in \sigma(X_k, X_{k+1}, \dots)$).

⁴ \mathcal{I} is the σ -algebra of *terminal events*;

Example 3.12. Let $(X_n)_{n \in \mathbb{N}_0}$ be independent random elements in a Polish space S (wlog defined on the sequence space), i.e.:

$$\mathbb{P} \equiv \mathbb{P}_X = \bigotimes_{n \in \mathbb{N}_0} \mathbb{P}_{X_n} .$$

Then we have $\mathbb{P}(A) \in \{0, 1\}$ for $A \in \mathcal{I}$; i.e. the shift $\varphi := \theta_1$ is ergodic.

Example 3.13 (Rotation on the circle). As in 3.4 we consider the transformation

$$\varphi : \Omega \longrightarrow \Omega, \quad \varphi(\omega) := \omega + \theta \pmod{1},$$

on the probability space $(\Omega, \mathcal{F}, \mathbb{P}) := ([0, 1), \mathcal{B}[0, 1), \lambda|_{\mathcal{F}})$, where λ denotes the Lebesgue measure. Then φ is ergodic iff θ is irrational.

PROOF. "‘ \Rightarrow ’" Let θ be rational, hence $\theta = \frac{m}{n}$ with integers $n \geq m \geq 1$. Moreover let $B \in \mathcal{F} \equiv \mathcal{B}[0, 1)$ with $0 < \lambda(B) < \frac{1}{n}$. Then $A := \bigcup_{k=1}^{m-1} (B + \frac{k}{n})$ is invariant, but $0 < \lambda(A) < 1$.

"‘ \Leftarrow ’" This can be proven with a Fourier series argument; see e.g. Shiryaev [Sh 95, p.408] or Kallenberg [KB 97, p.174/9]. \square

Example 3.14. Let $(X_n)_{n \in \mathbb{N}_0}$ be the canonical Markov chain on $S := \{1, 2, 3, 4\}$ with transition probability

$$p := \begin{pmatrix} \frac{1}{3} & \frac{2}{3} & 0 & 0 \\ \frac{2}{3} & \frac{1}{3} & 0 & 0 \\ 0 & 0 & \frac{1}{2} & \frac{1}{2} \\ 0 & 0 & \frac{1}{4} & \frac{3}{4} \end{pmatrix}$$

(p is a stochastic matrix, since the sums over the lines equal 1). A measure μ on S is invariant, if we have:

$$\mu(j) = \sum_{i=1}^4 p(i, j) \mu(i) \quad (j = 1, 2, 3, 4).$$

This is satisfied for instance by the measures

$$\mu_0(1) = \mu_0(2) := \frac{1}{2}, \quad \mu_0(3) = \mu_0(4) := 0$$

and

$$\mu_1(1) = \mu_1(2) := 0, \quad \mu_1(3) := \frac{1}{3}, \quad \mu_1(4) := \frac{2}{3}.$$

But then also each

$$\mu_\beta := (1 - \beta)\mu_0 + \beta\mu_1 \quad (0 \leq \beta \leq 1)$$

is invariant. With respect to the canonical shift $\varphi := \theta_1$ we have now:

$$A := \{X_n \in \{1, 2\}, n \in \mathbb{N}_0\} \in \mathcal{I} \quad \text{and} \quad B := \{X_n \in \{3, 4\}, n \in \mathbb{N}_0\} \in \mathcal{I}.$$

Hence we further have: $P_{\mu_\beta}(A) = 1 - \beta$ and $P_{\mu_\beta}(B) = \beta$. Consequently φ is ergodic iff $\beta \in \{0, 1\}$.

Theorem 3.15. *Let $(X_n)_{n \in \mathbb{N}_0}$ with Polish state space (S, \mathcal{S}) be ergodic and $g : S^{\mathbb{N}_0} \rightarrow S'$ be $\mathcal{S}^{\mathbb{N}_0}$ - \mathcal{S}' -measurable, where (S', \mathcal{S}') is equally Polish. Then*

$$Y_k := g(X_k, X_{k+1}, \dots) \quad (k \in \mathbb{N}_0)$$

is ergodic (in S').

PROOF. Wlog let again $(\Omega, \mathcal{F}, \mathbb{P}) = (S^{\mathbb{N}_0}, \mathcal{S}^{\mathbb{N}_0}, \mathbb{P}_{(X_n)_{n \in \mathbb{N}_0}})$ and $X_n = \pi_{\{n\}}$ (projection at time n) and $\varphi = \theta_1$ (shift). Equally let $(\Omega', \mathcal{F}', \mathbb{P}') = ((S')^{\mathbb{N}_0}, (\mathcal{S}')^{\mathbb{N}_0}, \mathbb{P}_{(Y_n)_{n \in \mathbb{N}_0}})$ and $\varphi' = \theta_1$. Moreover denote by \mathcal{I} resp. \mathcal{I}' the systems of invariant sets associated with φ resp. φ' .

Let now $A \in \mathcal{I}'$ be fixed; for $B := (g_0, g_1, \dots)^{-1}(A)$ we then have:

$$\begin{aligned} \varphi^{-1}(B) &= (g_1, g_2, \dots)^{-1}(A) \\ &= (g_0, g_1, \dots)^{-1}((\varphi')^{-1}(A)) \\ &= (g_0, g_1, \dots)^{-1}(A) \equiv B, \end{aligned}$$

hence $B \in \mathcal{I}$. By ergodicity of φ we get: $\mathbb{P}'(A) \equiv \mathbb{P}(B) \in \{0, 1\}$. \square

Example 3.16 (Bernoulli shift). As in 3.6 we consider i.i.d. random variables $(X_n)_n$ in $S := \{0, 1\}$ with $\mathbb{P}\{X_n = 0\} = \mathbb{P}\{X_n = 1\} = \frac{1}{2}$. Moreover let

$$g : \{0, 1\}^{\mathbb{N}_0} \rightarrow [0, 1), \quad (x_n)_n \mapsto \sum_{n=0}^{\infty} x_n 2^{-n-1} \pmod{1}.$$

By example 3.12 X is ergodic, hence according to Theorem 3.15 also

$$Y_k := g(X_k, X_{k+1}, \dots) \quad (k \in \mathbb{N}_0).$$

4. Birkhoff's ergodic theorem

Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space with measure preserving mapping $\varphi : \Omega \rightarrow \Omega$ and a random variable $X : \Omega \rightarrow \mathbb{R}$. We now study the asymptotic behavior of the stochastic process defined by

$$X_k := X \circ \varphi^k \quad (k \in \mathbb{N}_0).$$

Theorem 4.1 (Ergodic theorem, Birkhoff). *Let $X \in L^1(\mathbb{P})$. Then we have \mathbb{P} -a.s. and in $L^1(\mathbb{P})$:*

$$\frac{1}{n} \sum_{k=0}^{n-1} X \circ \varphi^k \xrightarrow{n \rightarrow \infty} \mathbb{E}(X | \mathcal{I}).$$

The proof is based on the following estimate:

Lemma 4.2 (Maximal-ergodic Lemma, Hopf). *In the situation from 4.1 let*

$$\begin{aligned} S_n &:= X_0 + \dots + X_{n-1} \equiv \sum_{k=0}^{n-1} X \circ \varphi^k && (n \in \mathbb{N}) \text{ and} \\ M_n &:= \max\{0, S_1, \dots, S_n\} && (n \in \mathbb{N}_0). \end{aligned}$$

Then we have:

$$\mathbb{E}(X 1_{\{M_n > 0\}}) \geq 0 \quad (n \in \mathbb{N}_0).$$

PROOF. In case $n = 0$ nothing is to be proven.

1. First we show:

$$X 1_{\{M_n > 0\}} \geq 1_{\{M_n > 0\}} (M_n - M_n \circ \varphi) \quad (n \in \mathbb{N});$$

by the definitions above we have $S_k - M_n \leq 0$ for all $k \in \{1, \dots, n\}$, hence also

$$X \geq X + (S_k - M_n) \circ \varphi = (X + S_k \circ \varphi) - M_n \circ \varphi \equiv S_{k+1} - M_n \circ \varphi$$

and thus

$$X \geq \max\{S_1, \dots, S_n\} - M_n \circ \varphi;$$

in particular we have shown:

$$\begin{aligned} X 1_{\{M_n > 0\}} &\geq 1_{\{M_n > 0\}} \max\{S_1, \dots, S_n\} - 1_{\{M_n > 0\}} M_n \circ \varphi \\ &= 1_{\{M_n > 0\}} (M_n - M_n \circ \varphi) \end{aligned} \quad (n \in \mathbb{N}),$$

hence the desired claim.

2. But this implies:

$$\begin{aligned} \mathbb{E}(X 1_{\{M_n > 0\}}) &\geq \int_{\{M_n > 0\}} (M_n - M_n \circ \varphi) d\mathbb{P} \\ &= \int (M_n - M_n \circ \varphi) d\mathbb{P} = 0, \end{aligned}$$

where we finally use that φ is measure preserving. □

PROOF OF BIRKHOFF'S ERGODIC THEOREM 4.1 Wlog assume $\mathbb{E}(X|\mathcal{I}) = 0$; else we consider $\tilde{X} := X - \mathbb{E}(X|\mathcal{I})$, which is possible thanks to the invariance of $\mathbb{E}(X|\mathcal{I}) \circ \varphi = \mathbb{E}(X|\mathcal{I})$ (\mathbb{P} -a.s.).

1. \mathbb{P} -almost sure convergence: For this purpose we will show that with

$$\bar{X} := \limsup_{n \rightarrow \infty} \frac{S_n}{n} \equiv \limsup_{n \rightarrow \infty} \frac{1}{n} \sum_{k=0}^{n-1} X \circ \varphi^k$$

and with

$$D := \{ \bar{X} > \varepsilon \} \in \mathcal{I} \quad (\text{for arbitrary } \varepsilon > 0)$$

we have

$$\mathbb{P}(D) = 0;$$

analogously one can show $\liminf \frac{S_n}{n} \geq 0$, by considering $-X$ instead of X . To show $\mathbb{P}(D) = 0$ we give D an alternative description: With

$$\begin{aligned} X^* &:= (X - \varepsilon) 1_D \\ S_n^* &:= X^* + X^* \circ \varphi + \dots + X^* \circ \varphi^{n-1} \\ M_n^* &:= \max\{0, S_1^*, \dots, S_n^*\} \\ F_n &:= \{M_n^* > 0\} \end{aligned}$$

we have

$$D = \left\{ \sup_{n \in \mathbb{N}} \frac{S_n^*}{n} > 0 \right\} = \bigcup_{n \in \mathbb{N}} F_n.$$

Upon applying the maximal-ergodic Lemma 4.2 on X^* , we obtain:

$$\begin{aligned} 0 &\leq \mathbb{E}(X^* 1_{F_n}) && (\text{Lemma 4.2}) \\ &\xrightarrow{n \rightarrow \infty} \mathbb{E}\left(X^* 1_{\bigcup_n F_n}\right) && (\text{dom. conv., since } X \in L^1) \\ &= \mathbb{E}(X^* 1_D) && (\text{above char. of } D) \\ &\equiv \mathbb{E}(X 1_D) - \varepsilon \mathbb{P}(D) && (\text{definition of } X^*) \\ &= -\varepsilon \mathbb{P}(D) && (\mathbb{E}(X|\mathcal{I}) = 0 \text{ and } D \in \mathcal{I}) \\ &\leq 0, \end{aligned}$$

in summary: $\mathbb{P}(D) = 0$.

2. L^1 -convergence: For this purpose we truncate X ; with some fixed $K > 0$ let

$$X' := X 1_{\{|X| \leq K\}} \quad \text{and} \quad X'' := X - X'.$$

The \mathbb{P} -a.s. convergence proved above applies in particular to X' ; since this convergence is dominated by K , we obtain for X' :

$$\frac{1}{n} \sum_{k=0}^{n-1} X' \circ \varphi^k \xrightarrow{n \rightarrow \infty} \mathbb{E}(X'|\mathcal{I}) \quad \text{in } L^1(\mathbb{P}).$$

In addition we have

$$\mathbb{E} \left(\left| \frac{1}{n} \sum_{k=0}^{n-1} X'' \circ \varphi^k \right| \right) \leq \frac{1}{n} \sum_{k=0}^{n-1} \mathbb{E} (|X'' \circ \varphi^k|) = \mathbb{E}(|X''|),$$

where we use that φ is measure preserving; moreover we have by Jensen ($|\cdot|$ is convex):

$$\mathbb{E} (|\mathbb{E}(X''|\mathcal{F})|) \leq \mathbb{E} (\mathbb{E}(|X''| |\mathcal{F})) = \mathbb{E}(|X''|);$$

combining the preceding two inequalities we obtain

$$\mathbb{E} \left(\left| \frac{1}{n} \sum_{k=0}^{n-1} X'' \circ \varphi^k - \mathbb{E}(X''|\mathcal{F}) \right| \right) \leq 2 \mathbb{E}(|X''|).$$

Let now an arbitrary $\varepsilon > 0$ be fixed; then we may choose $K > 0$ such that

$$2 \mathbb{E}(|X''|) < \frac{\varepsilon}{2}$$

(dominated convergence, definition of X''). With the parameters ε and K by the above L^1 -convergence for X' we can choose $n_0 \in \mathbb{N}$, such that we have:

$$\mathbb{E} \left(\left| \frac{1}{n} \sum_{k=0}^{n-1} X' \circ \varphi^k - \mathbb{E}(X'|\mathcal{F}) \right| \right) < \frac{\varepsilon}{2} \quad (n \geq n_0).$$

Since now $X \equiv X' + X''$, the preceding three estimates yield:

$$\begin{aligned} & \mathbb{E} \left(\left| \frac{1}{n} \sum_{k=0}^{n-1} X \circ \varphi^k - \mathbb{E}(X|\mathcal{F}) \right| \right) \\ & \leq \mathbb{E} \left(\left| \frac{1}{n} \sum_{k=0}^{n-1} X' \circ \varphi^k - \mathbb{E}(X'|\mathcal{F}) \right| \right) + \mathbb{E} \left(\left| \frac{1}{n} \sum_{k=0}^{n-1} X'' \circ \varphi^k - \mathbb{E}(X''|\mathcal{F}) \right| \right) < \varepsilon. \end{aligned}$$

□

Example 4.3 (Strong law of large numbers). Let $(X_n)_{n \in \mathbb{N}_0}$ be i.i.d. random variables, wlog defined on the sequence space $\Omega := \mathbb{R}^{\mathbb{N}_0}$, with $\mathbb{P} \equiv P_X = P_{X_0} \otimes P_{X_0} \otimes \cdots$ and ergodic shift $\varphi = \theta_1$ (see example 3.12). If $X_0 \in L^1(\mathbb{P})$, we infer from Theorem 4.1 with Proposition 3.9:

$$\frac{1}{n} \sum_{k=0}^{n-1} X_k = \frac{1}{n} \sum_{k=0}^{n-1} X_0 \circ \varphi^k \xrightarrow{\mathbb{P}\text{-a.s., } L^1(\mathbb{P})} \mathbb{E}(X_0|\mathcal{F}) = \mathbb{E}(X_0).$$

Example 4.4 (Rotation on the circle, Weyl's equidistribution law). Let

$$\varphi : \Omega \longrightarrow \Omega, \quad \varphi(\omega) := \omega + \theta \pmod{1},$$

on $(\Omega, \mathcal{F}, \mathbb{P}) := ([0, 1), \mathcal{B}[0, 1), \lambda|_{\mathcal{F}})$ be given as in 3.4 and 3.13, where λ is the Lebesgue measure. Moreover let $\theta \in \mathbb{Q}^c$. Then from Theorem 4.1 with example 3.13 for $A \in \mathcal{B}[0, 1)$

we obtain:

$$\frac{1}{n} \sum_{k=0}^{n-1} 1_A \circ \varphi^k \xrightarrow{\lambda\text{-a.e., } L^1(\lambda)} \lambda(A).$$

5. The subadditive ergodic theorem of Kingman

Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space with a measure preserving transformation $\varphi : \Omega \rightarrow \Omega$. In the preceding section we have investigated the asymptotic behavior of $\frac{S_n}{n}$, where S_n is of the form $\sum_{k=0}^{n-1} X \circ \varphi^k$, so in particular satisfies the *additive cocycle property*

$$S_{n+m} = S_n + S_m \circ \varphi^n \quad (n, m \in \mathbb{N}_0).$$

Now we are interested in the following generalization:

Definition 5.1 (subadditive sequences of random variables). A sequence $(Y_n)_n$ of random variables ($n \in \mathbb{N}_0$ or \mathbb{N} ; state space $\mathbb{R} \cup \{-\infty\}$) is called *subadditive*, if we have:

$$Y_{n+m} \leq Y_n + Y_m \circ \varphi^n \quad (n, m \in \mathbb{N}_0).$$

A sequence $(Y_n)_{n \in \mathbb{N}_0}$ is called *superadditive*, if $(-Y_n)_{n \in \mathbb{N}_0}$ is subadditive, and it is called *additive*, if it is both sub- and superadditive.

Example 5.2. Let $(X_n)_{n \in \mathbb{N}_0}$ be a sequence of i.i.d. random variables, realized wlog by $X_n = \pi_{\{n\}}$ on the sequence space $(\Omega, \mathcal{F}, \mathbb{P}, \varphi) = (S^{\mathbb{N}_0}, \mathcal{S}^{\mathbb{N}_0}, P_{(X_n)_{n \in \mathbb{N}_0}}, \theta_1)$. Let

$$S_n := \sum_{k=0}^{n-1} X_k.$$

Then $(S_n)_{n \in \mathbb{N}_0}$ is additive and $(|S_n|)_{n \in \mathbb{N}_0}$ subadditive.

PROOF. Additivity of $(S_n)_n$ is immediate, since $\varphi \equiv \theta_1$. Moreover we have:

$$\begin{aligned} |S_{n+m}| &\equiv \left| \sum_{k=0}^{n+m-1} X_k \right| \leq \left| \sum_{k=0}^{n-1} X_k \right| + \left| \sum_{k=n}^{n+m-1} X_k \right| \\ &= |S_n| + \left| \sum_{k=0}^{m-1} X_k \circ \varphi^n \right| = |S_n| + |S_m| \circ \varphi^n. \end{aligned}$$

□

Example 5.3. (discrete version of linear stochastic differential equation) Let B_0, B_1 be $d \times d$ matrices with real values, W a one-dimensional Brownian motion, θ_1 the shift by time one on Wiener space. We consider the discrete version

$$x_{n+1} - x_n = B_0 x_n + B_1 (W_{n+1} - W_n)$$

of the stochastic differential equation

$$dx_t = B_0 x_t dt + B_1 x_t dW_t.$$

The discrete equation may be written

$$x_{n+1} = (I + B_0 + B_1(W_1 - W_0) \circ \theta^n) x_n = (A \circ \theta^n) x_n$$

with the random matrix $A : \Omega \rightarrow \mathbb{R}^{d \times d}$ given by $A = (I + B_0 + B_1(W_1 - W_0))$. This is a special case of the following example. On a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ with measure preserving map φ let a random matrix be given, i.e. a measurable mapping $A : \Omega \rightarrow \mathbb{R}^{d \times d}$. Moreover let

$$\begin{aligned} A_n &:= (A \circ \varphi^{n-1})(A \circ \varphi^{n-2}) \cdots A && \text{and thus} \\ Y_n &:= \log \|A_n\| && (n \in \mathbb{N}), \end{aligned}$$

where $\|\cdot\|$ denotes a matrix norm. Then $(Y_n)_n$ is subadditive.

PROOF.

$$\begin{aligned} Y_{n+m} &= \log \|(A \circ \varphi^{m-1} \circ \varphi^n) \cdots (A \circ \varphi^0 \circ \varphi^n)(A \circ \varphi^{n-1}) \cdots A\| \\ &= \log \|(A_m \circ \varphi^n) A_n\| \\ &\stackrel{\text{Norm}}{\leq} \log [(\|A_m\| \circ \varphi^n) \|A_n\|] \\ &= \log (\|A_m\| \circ \varphi^n) + \log \|A_n\| \equiv Y_m \circ \varphi^n + Y_n. \end{aligned}$$

□

Now we aim at getting a convergence statement for subadditive $(Y_n)_n$. This will be achieved in the subadditive ergodic theorem 5.7 of Kingman. For this we state the three following Lemmas.

Lemma 5.4 (Riesz). *Let $u_1, \dots, u_n \in \mathbb{R}$ ($n \in \mathbb{N}$). With*

$$s_j := \begin{cases} 0, & j = 0 \\ u_1 + \cdots + u_j, & j \in \{1, \dots, n\}, \end{cases}$$

define

$$v_j \equiv v_{jn} := \max_{k \in \{j, \dots, n\}} (s_k - s_j) \equiv \max\{0, u_{j+1}, u_{j+1} + u_{j+2}, u_{j+1} + \cdots + u_n\}$$

for $j = 0, 1, \dots, n$. Then we have:

$$\sum_{j=0}^{n-1} u_{j+1} 1_{\{v_{jn} > 0\}} \geq 0.$$

PROOF. 1) First we have for all $j \in \{0, 1, \dots, n\}$:

$$v_j = \max\{0, u_{j+1} + v_{j+1}\} \equiv (u_{j+1} + v_{j+1})^+.$$

This follows directly, since

$$\begin{aligned} v_j &= \max\{0, u_{j+1}, u_{j+1} + u_{j+2}, u_{j+1} + \dots + u_n\} \quad \text{and} \\ &= \max\{0, \max\{u_{j+1}, u_{j+1} + u_{j+2}, \dots, u_{j+1} + \dots + u_n\}\} \\ &= \max\{0, u_{j+1} + \max\{0, u_{j+2}, u_{j+2} + u_{j+3}, \dots, u_{j+2} + \dots + u_n\}\} \\ &= \max\{0, u_{j+1} + v_{j+1}\}. \end{aligned}$$

2) By 1) we have:

$$v_j \leq v_{j+1} + u_{j+1} \mathbf{1}_{\{v_j > 0\}} \quad (j \in \{0, 1, \dots, n\}).$$

Indeed, if $v_j = 0$, this is trivial, and in case $v_j > 0$ we have:

$$0 < v_j \stackrel{1)}{=} (u_{j+1} + v_{j+1})^+ \stackrel{v_j > 0}{=} v_{j+1} + u_{j+1}.$$

3) From 2) now follows the claim of the Lemma, since:

$$0 \leq v_0 = v_0 - v_n = \sum_{j=0}^{n-1} (v_j - v_{j+1}) \stackrel{2)}{\leq} \sum_{j=0}^{n-1} u_{j+1} \mathbf{1}_{\{v_j > 0\}}.$$

□

In the proof of the subadditive ergodic theorem of Kingman we will compare subadditive sequences $(Y_n)_n$ with additive sequences $X_n = \sum_{i=0}^{n-1} X_0 \circ \varphi^i$. For this purpose we need the following auxiliary argument, for which the preceding Lemma of Riesz will be useful.

Lemma 5.5 (maximal inequality). *Let $(Y_n)_{n \in \mathbb{N}_0}$ be superadditive on $(\Omega, \mathcal{F}, \mathbb{P}, \varphi)$ and assume $Y_n \geq 0$ for all n . Moreover let $X \geq 0$ be an integrable random variable; set*

$$V := \sup_{n \in \mathbb{N}_0} (Y_n - X_n) - Y_0, \quad \text{where} \quad X_n := \sum_{i=0}^{n-1} X \circ \varphi^i.$$

Then we have:

$$\mathbb{E}(X \mathbf{1}_{\{V > 0\}} | \mathcal{F}) \leq \sup_{n \in \mathbb{N}} \frac{\mathbb{E}(Y_n | \mathcal{F})}{n}.$$

PROOF. Let $v_{jn} := \max_{k \in \{j, \dots, n\}} \left(Y_k - Y_j - \sum_{i=j}^{k-1} X \circ \varphi^i \right)$ for $j = 0, 1, \dots, n$.

1) At first we have:

$$Y_n \geq \sum_{j=0}^{n-1} X \circ \varphi^j \mathbf{1}_{\{v_{jn} > 0\}} \quad (n \in \mathbb{N});$$

since with $Y_{j+1} \geq Y_j$ (by superadditivity and $Y_n \geq 0$) we obtain:

$$\begin{aligned} Y_n &\geq Y_n - Y_0 = \sum_{j=0}^{n-1} (Y_{j+1} - Y_j) \\ &\geq \sum_{j=0}^{n-1} (Y_{j+1} - Y_j) \mathbf{1}_{\{v_{jn} > 0\}} \stackrel{5.4}{\geq} \sum_{j=0}^{n-1} X \circ \varphi^j \mathbf{1}_{\{v_{jn} > 0\}}, \end{aligned}$$

where the last step follows from Lemma 5.4 with $u_j := Y_j - Y_{j-1} - X \circ \varphi^{j-1}$.

2) From 1) we get:

$$\mathbb{E}(Y_n | \mathcal{F}) \geq \sum_{k=1}^n \mathbb{E}(X \mathbf{1}_{\{v_{0k} > 0\}} | \mathcal{F}) \quad (n \in \mathbb{N});$$

since for $k \geq j$ by superadditivity $Y_k - Y_j \geq Y_{k-j} \circ \varphi^j$ and thus $v_{jn} \geq v_{0(n-j)} \circ \varphi^j$, hence in summary (with the measure preserving property of φ):

$$\begin{aligned} \mathbb{E}(Y_n | \mathcal{F}) &\stackrel{1)}{\geq} \sum_{j=0}^{n-1} \mathbb{E}\left(X \circ \varphi^j \mathbf{1}_{\{v_{jn} > 0\}} \mid \mathcal{F}\right) \\ &\geq \sum_{j=0}^{n-1} \mathbb{E}\left(\left[X \mathbf{1}_{\{v_{0(n-j)} > 0\}}\right] \circ \varphi^j \mid \mathcal{F}\right) = \sum_{k=1}^n \mathbb{E}(X \mathbf{1}_{\{v_{0k} > 0\}} | \mathcal{F}). \end{aligned}$$

3) By Fatou and $\{v_{0k} > 0\} \nearrow \{V > 0\}$ ($k \rightarrow \infty$) we obtain from this:

$$\sup_{n \in \mathbb{N}} \frac{\mathbb{E}(Y_n | \mathcal{F})}{n} \stackrel{2)}{\geq} \liminf_{n \rightarrow \infty} \frac{1}{n} \sum_{k=1}^n \mathbb{E}(X \mathbf{1}_{\{v_{0k} > 0\}} | \mathcal{F}) \geq \mathbb{E}(X \mathbf{1}_{\{V > 0\}} | \mathcal{F}).$$

□

Lemma 5.6. *If Z is a measurable function on $(\Omega, \mathcal{F}, \mathbb{P}, \varphi)$ with $Z \geq Z \circ \varphi$, we have $Z = Z \circ \varphi$. If in particular $(Y_n)_{n \in \mathbb{N}_0}$ is superadditive and*

$$\bar{Y} := \limsup_{n \rightarrow \infty} \frac{Y_n}{n}, \quad \text{resp.} \quad \underline{Y} := \liminf_{n \rightarrow \infty} \frac{Y_n}{n},$$

we have:

$$\bar{Y} = \bar{Y} \circ \varphi, \quad \text{resp.} \quad \underline{Y} = \underline{Y} \circ \varphi.$$

PROOF. We first prove the statement on Z and for this assume $Z > Z \circ \varphi$ on a set of positive measure, so

$$\mathbb{P}(Z > q > Z \circ \varphi) > 0$$

for some $q \in \mathbb{Q}$. Then we get:

$$\begin{aligned} \mathbb{P}(Z < q) &\stackrel{\varphi \text{ m.p.}}{=} \mathbb{P}(Z \circ \varphi < q) \\ &= \mathbb{P}(Z \circ \varphi < q \leq Z) + \mathbb{P}(Z \circ \varphi < q, Z < q) \\ &\stackrel{Z \geq Z \circ \varphi}{=} \underbrace{\mathbb{P}(Z \circ \varphi < q \leq Z)}_{>0} + \mathbb{P}(Z < q) \\ &> \mathbb{P}(Z < q), \end{aligned}$$

a contradiction.

With what has just been proved, it now remains to show:

$$\bar{Y} \geq \bar{Y} \circ \varphi, \quad \text{resp.} \quad \underline{Y} \geq \underline{Y} \circ \varphi;$$

but by the superadditivity of $(Y_n)_n$ we now have:

$$\begin{aligned} \frac{Y_{n+1}}{n+1} &\geq \frac{Y_1}{n+1} + \frac{Y_n \circ \varphi}{n+1} \\ &= \frac{Y_1}{n+1} + \frac{n}{n+1} \frac{Y_n}{n} \circ \varphi. \end{aligned}$$

□

Theorem 5.7 (subadditive ergodic theorem, Kingman).

On $(\Omega, \mathcal{F}, \mathbb{P}, \varphi)$ let $(Y_n)_{n \in \mathbb{N}}$ be a superadditive sequence of integrable random variables. Then we have:

$$\frac{Y_n}{n} \xrightarrow[n \rightarrow \infty]{\mathbb{P}\text{-a.s.}} \sup_{n \in \mathbb{N}} \frac{1}{n} \mathbb{E}(Y_n | \mathcal{I}) =: \gamma \leq \infty.$$

Here γ is integrable iff $\sup_{n \in \mathbb{N}} \frac{1}{n} \mathbb{E}(Y_n) < \infty$. In this case we also have

$$\frac{Y_n}{n} \xrightarrow[n \rightarrow \infty]{L^1(\mathbb{P})} \gamma.$$

Moreover there exists a set $\tilde{\Omega} \in \mathcal{I}$ with $\tilde{\Omega} \subset \varphi^{-1} \tilde{\Omega}$ and $\mathbb{P}(\tilde{\Omega}) = 1$, so that we also have:

$$\frac{Y_n}{n} \xrightarrow[n \rightarrow \infty]{} \gamma \quad \text{on } \tilde{\Omega}.$$

PROOF. To simplify notation we set $Y_0 := 0$ (else replace Y_n by $Y_n - Y_0$, $n \geq 0$); then $(Y_n)_{n \in \mathbb{N}_0}$ is furthermore superadditive.

1) We first show that we can assume wlog

$$Y_n \geq 0 \quad (n \in \mathbb{N}).$$

For this purpose let

$$G_n := Y_n - F_n, \quad \text{with} \quad F_n := \sum_{i=0}^{n-1} Y_1 \circ \varphi^i.$$

Then

$$Y_n = G_n + F_n.$$

The sequence $(F_n)_n$ is an additive cocycle, so that the assumed integrability provides due to Birkhoff's ergodic theorem 4.1 a random variable γ_F such that

$$\frac{F_n}{n} \longrightarrow \gamma_F$$

\mathbb{P} -a.s. and in $L^1(\mathbb{P})$. The claims about Y follow, if we have shown

$$\frac{G_n}{n} \longrightarrow \gamma_G \quad (\mathbb{P}\text{-a.s. and in } L^1(\mathbb{P})),$$

since then in particular

$$\frac{Y_n}{n} \longrightarrow \gamma_G + \gamma_F \quad (\mathbb{P}\text{-a.s. and in } L^1(\mathbb{P})).$$

But now by inductive application of superadditivity

$$\begin{aligned} G_n &= -Y_1 - Y_1 \circ \varphi - \dots - Y_1 \circ \varphi^{n-2} - Y_1 \circ \varphi^{n-1} + Y_n \\ &\geq -Y_1 - Y_1 \circ \varphi - \dots - Y_1 \circ \varphi^{n-2} + Y_{n-1} \\ &\geq -Y_1 - Y_1 \circ \varphi - \dots + Y_{n-2} \\ &\geq \dots \dots \\ &\geq -Y_1 - Y_1 \circ \varphi + Y_2 \geq 0; \end{aligned}$$

on the other hand superadditivity of $(Y_n)_n$ is transferred to $(G_n)_n$. Therefore the claims for $(Y_n)_n$ reduce to corresponding convergences for the positive process $(G_n)_n$.

2) We further prove that wlog we can assume

$$Y_n \geq n \quad (n \in \mathbb{N}) :$$

Since $Y_{n+m} + n + m \geq Y_n + n + (Y_m + m) \circ \varphi^n$ with $(Y_n)_n$ also $(Y_n + n)_n$ is superadditive. If $\frac{Y_n + n}{n} \longrightarrow \gamma'$, so also $\frac{Y_n}{n} \longrightarrow \gamma := \gamma' - 1$.

3) $\frac{Y_n}{n} \rightarrow \gamma$ \mathbb{P} -a.s.: For this we show $\bar{Y} \leq \gamma$ and $\underline{Y} \geq \gamma$, where again

$$\bar{Y} := \limsup_{n \rightarrow \infty} \frac{Y_n}{n} \quad \text{resp.} \quad \underline{Y} := \liminf_{n \rightarrow \infty} \frac{Y_n}{n}.$$

i) $\bar{Y} \leq \gamma$ \mathbb{P} -a.s.: For $2 < r \in \mathbb{N}$ we define

$$X^r := \min\left\{r, \bar{Y} - \frac{1}{r}\right\} > 0;$$

with this the inequality follows, since by 2) $\bar{Y} \geq 1$. By Lemma 5.6 we have moreover $X^r = X^r \circ \varphi$, hence $X_n^r := \sum_{i=0}^{n-1} X^r \circ \varphi^i = nX^r$; with this we have

$$V := \sup_{n \in \mathbb{N}_0} (Y_n - X_n^r) - Y_0 = \sup_{n \in \mathbb{N}_0} (Y_n - nX^r) > 0;$$

here the latter inequality follows from the definition of X^r : in case $Y_n \leq nX^r$ for all $n \in \mathbb{N}_0$, the contradiction $\bar{Y} \equiv \limsup \frac{Y_n}{n} \leq \limsup \frac{nX^r}{n} = X^r < \bar{Y}$ would follow. Hence with Lemma 5.5:

$$X^r = \mathbb{E}(X^r | \mathcal{I}) \leq \sup_{n \in \mathbb{N}} \frac{\mathbb{E}(Y_n | \mathcal{I})}{n} \equiv \gamma.$$

But from this as $r \rightarrow \infty$: $\bar{Y} \leq \gamma$ \mathbb{P} -a.s..

ii) $\underline{Y} \geq \gamma$ \mathbb{P} -a.s.: First we have that $(Y_n)_n$ is increasing, by superadditivity and positivity; from this we conclude:

$$k Y_{n+k-1} \geq \sum_{l=1}^k \sum_{j=0}^{n-1} (Y_{j+l} - Y_{j+l-1}) = \sum_{j=0}^{n-1} (Y_{j+k} - Y_j) \quad (k, n \in \mathbb{N});$$

consequently for each $k \in \mathbb{N}$:

$$\begin{aligned} \underline{Y} &= \liminf_{n \rightarrow \infty} \frac{Y_{n+k-1}}{n+k-1} \\ &= \liminf_{n \rightarrow \infty} \frac{Y_{n+k-1}}{n} \\ &= \frac{1}{k} \liminf_{n \rightarrow \infty} \frac{k Y_{n+k-1}}{n} \\ &\geq \frac{1}{k} \liminf_{n \rightarrow \infty} \sum_{j=0}^{n-1} \frac{Y_{j+k} - Y_j}{n} && \text{(prec. rem.)} \\ &\geq \frac{1}{k} \liminf_{n \rightarrow \infty} \sum_{j=0}^{n-1} \frac{Y_k \circ \varphi^j}{n} && \text{(superadditivity)} \\ &= \frac{1}{k} \mathbb{E}(Y_k | \mathcal{I}) && \text{(Birkhoff 4.1)} \end{aligned}$$

hence also $\underline{Y} \geq \gamma$.

4) γ integrable $\Leftrightarrow \sup_{n \in \mathbb{N}} \frac{1}{n} \mathbb{E}(Y_n) < \infty$: For this let $Z_n := \frac{Y_n}{n}$. From what has been shown: $Z_n \rightarrow \gamma$ \mathbb{P} -a.s. and $\mathbb{E}\gamma \geq \mathbb{E}Z_n$; therefore " \Rightarrow " has been shown. " \Leftarrow " follows with monotone convergence.

5) If γ is integrable, we have: $Z_n \equiv \frac{Y_n}{n} \rightarrow \gamma$ in $L^1(\mathbb{P})$. In fact, by $0 \leq (\gamma - Z_n)^+ \leq \gamma$ we have on the one hand

$$\mathbb{E}((\gamma - Z_n)^+) \rightarrow 0;$$

and on the other hand also

$$0 \leq \mathbb{E}(\gamma - Z_n) \rightarrow 0$$

by Fatou, so that we also obtain:

$$\mathbb{E}((\gamma - Z_n)^-) = -\mathbb{E}(\gamma - Z_n) + \mathbb{E}((\gamma - Z_n)^+) \rightarrow 0,$$

and thus in summary

$$\mathbb{E}(|\gamma - Z_n|) \rightarrow 0.$$

6) existence of $\tilde{\Omega} \in \mathcal{S}$ with $\tilde{\Omega} \subset \varphi^{-1}\tilde{\Omega}$, $\mathbb{P}(\tilde{\Omega}) = 1$ and $\frac{Y_n}{n} \rightarrow \gamma$ on $\tilde{\Omega}$: By Lemma 5.6 \bar{Y} and \underline{Y} are invariant. Therefore also

$$\tilde{\Omega} := \{\bar{Y} = \underline{Y}\}$$

is invariant; the remaining properties follow from what has been proved (evtl. choice of an a.s. equal strictly invariant set). \square

6. The theorem of Furstenberg-Kesten

Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space with a measure preserving mapping $\varphi : \Omega \rightarrow \Omega$ and $A : \Omega \rightarrow \mathbb{R}^{d \times d}$ a random matrix. We now investigate the asymptotics of

$$A_n(\omega) := (A \circ \varphi^{n-1}(\omega)) (A \circ \varphi^{n-2}(\omega)) \cdots (A \circ \varphi(\omega)) (A(\omega)) \quad (\omega \in \Omega). \quad (2)$$

Example 6.1 (deterministic, symmetric matrix). Let $A \in \mathbb{R}^{d \times d}$ be symmetric. Then there exists a diagonalisation of A with real (A symmetric) eigenvalues $\delta_1 \geq \cdots \geq \delta_d$; hence there exists an orthogonal matrix O , so that we have

$$A = O^* D O \quad \text{with} \quad D := \begin{pmatrix} \delta_1 & & 0 \\ & \ddots & \\ 0 & & \delta_d \end{pmatrix}.$$

Here we assume $\delta_1 > \cdots > \delta_d$, i.e. the eigenspaces E_i belonging to δ_i are one dimensional. Moreover let x_i be an eigenvector in E_i and

$$V_j = \begin{cases} E_j \oplus E_{j+1} \oplus \cdots \oplus E_d, & j = 1, \dots, d, \\ \{0\}, & j = d+1. \end{cases}$$

In this setting let $x \in V_j \setminus V_{j+1}$. Then x can be written as

$$x = \sum_{k=j}^d \alpha_k x_k \quad \text{with} \quad \alpha_j \neq 0.$$

So we have by linearity of A , since x_k are eigenvectors:

$$A^n x = \sum_{k=j}^d \alpha_k A^n x_k = \sum_{k=j}^d \alpha_k \delta_k^n x_k,$$

hence

$$\begin{aligned} \frac{1}{n} \log |A^n x| &= \frac{1}{n} \log \left| \sum_{k=j}^d \alpha_k \delta_k^n x_k \right| \\ &= \frac{1}{n} \left[\log \delta_j^n + \log \left| \sum_{k=j}^d \alpha_k \left(\frac{\delta_k}{\delta_j} \right)^n x_k \right| \right] \xrightarrow{n \rightarrow \infty} \log \delta_j. \end{aligned}$$

We also have the reverse conclusion, in summary:

$$x \in V_j \setminus V_{j+1} \iff \lim_{n \rightarrow \infty} \frac{1}{n} \log \|A^n x\| = \log \delta_j \quad (j = 1, \dots, d).$$

We now aim at showing this analogously also for the sequence $(A^n)_n$ from (2).

Definition-remark 6.2 (decomposition according to singular value). Each $A \in \mathbb{R}^{d \times d}$ possesses a *singular value decomposition*, i.e. there are orthogonal matrices U, V and a diagonal matrix

$$D = \begin{pmatrix} \delta_1 & & 0 \\ & \ddots & \\ 0 & & \delta_d \end{pmatrix}$$

with $\delta_1 \geq \dots \geq \delta_d$, such that we have

$$A = VDU.$$

Thereby $\delta_1, \dots, \delta_d$ are the eigenvalues of $(A^*A)^{1/2}$ and for the operator norm we have: $\|A\| = \delta_1$.

PROOF. First of all, A possesses a *polar decomposition*, i.e.:

$$A = W(A^*A)^{1/2} \quad \text{with an orthogonal matrix } W.$$

(In case A is non-singular, this follows with $W := A(A^*A)^{-1/2}$). Let now $D := \text{diag}(\delta_1, \dots, \delta_d)$ be the diagonal matrix with eigenvalues $\delta_1 \geq \dots \geq \delta_d$ of $(A^*A)^{1/2}$. So the positive semidefinite matrix $(A^*A)^{1/2}$ can be written as

$$(A^*A)^{1/2} = U^*DU$$

with an orthogonal matrix U . In summary we have with $V := WU^*$

$$A = W(A^*A)^{1/2} = WU^*DU = VDU.$$

□

Remark: If e_1, \dots, e_d are the canonical unit vectors in \mathbb{R}^d , then Ue_i is the vector in direction of the i -th main axis of the ellipsoid $(A^*A)^{1/2}(S^{d-1})$ and δ_i describes the dilation in this direction.

For the construction of the analogues of $\delta_1, \dots, \delta_d$ in example 6.1 for the sequence A_n defined in (2) we need information on how A_n acts on linear subspaces of \mathbb{R}^d of any dimension below d . These linear subspaces are elements of the Grassmannian manifolds which can be defined via exterior products. In the Theorem by Furstenberg-Kesten we shall investigate the asymptotic behavior of the singular values $\delta_i(A_n)$ for $1 \leq i \leq d$ as $n \rightarrow \infty$.

Definition 6.3 (exterior product). For a d -dimensional linear space E let E^* be the dual of E and

$$\mathcal{L}^k(E^*) := \{k\text{-linear forms on } (E^*)^k\} \quad (k = 1, \dots, d).$$

We thus define $\wedge^k E$, the k -fold "exterior product of E ", as

$$\wedge^k E := \{f \in \mathcal{L}^k(E^*) : f \text{ alternating}\},$$

hence as the collection of all k -linear, alternating multilinear forms on $(E^*)^k$.

An element $f \in \wedge^k E$ is a k -linear mapping

$$f : \underbrace{E^* \times \dots \times E^*}_{k \text{ times}} \longrightarrow \mathbb{R},$$

which is alternating, i.e.:

$$f(\dots, x_i, \dots, x_j, \dots) = -f(\dots, x_j, \dots, x_i, \dots) \quad (i \neq j).$$

Lemma 6.4 (alternating maps). *For $f \in \mathcal{L}^k(E^*)$ the following are equivalent:*

- i) $f \in \wedge^k E$;
- ii) $f(x_1, \dots, x_k) = 0$, if (x_1, \dots, x_k) not pairwise different;
- iii) $f(x_1, \dots, x_k) = 0$, if (x_1, \dots, x_k) not pairwise linearly independent;
- iv) $f(x_{\pi(1)}, \dots, x_{\pi(k)}) = \text{sgn}(\pi) f(x_1, \dots, x_k)$ for all $\pi \in \mathfrak{S}_k$.

PROOF. i) \Leftrightarrow iv) follows from the representation $\pi = \tau_1 \circ \dots \circ \tau_k$ with permutations of two elements τ_i , hence $\text{sgn}(\tau_i) = -1$.

i) \Rightarrow ii): If (x_1, \dots, x_k) are not pairwise different, from the definition of an alternating mapping we obtain by exchanging equal elements:
 $f(x_1, \dots, x_k) = -f(x_1, \dots, x_k)$.

ii) \Rightarrow iii): Let wlog $x_k = \sum_{i=1}^{k-1} \alpha_i x_i$. Then with linearity and ii):
 $f(x_1, \dots, x_k) = \sum_{i=1}^{k-1} \alpha_i f(x_1, \dots, x_{k-1}, x_i) = 0$.

iii) \Rightarrow ii) is trivial.

ii) \Rightarrow i): Let wlog $k = 2$. Then we have for $x_1, x_2 \in E^*$:

$$\begin{aligned} 0 &\stackrel{\text{ii)}}{=} f(x_1 + x_2, x_1 + x_2) \\ &= f(x_1, x_1) + f(x_1, x_2) + f(x_2, x_1) + f(x_2, x_2) \\ &\stackrel{\text{ii)}}{=} f(x_1, x_2) + f(x_2, x_1), \end{aligned}$$

hence $f(x_1, x_2) = -f(x_2, x_1)$. □

Definition-remark 6.5. Let $f \in \wedge^k E$ and $g \in \wedge^l E$, where E denotes again a d -dimensional linear space and $k, l \in \mathbb{N}_0$. Then

$$f \wedge g(x_1, \dots, x_{k+l}) := \frac{1}{k!l!} \sum_{\pi \in \mathfrak{S}_{k+l}} \text{sgn}(\pi) f(x_{\pi(1)}, \dots, x_{\pi(k)}) g(x_{\pi(k+1)}, \dots, x_{\pi(k+l)})$$

is called the *exterior product of f and g* and we have: $f \wedge g \in \wedge^{k+l} E$.

PROOF. We have: $f \wedge g \in \mathcal{L}^m(E^*)$ with $m := k + l$. To see that $f \wedge g$ is alternating, we apply 6.4 iv); for arbitrary $x_1, \dots, x_m \in E^*$ let

$$a(\pi) := \frac{1}{k!l!} f(x_{\pi(1)}, \dots, x_{\pi(k)}) g(x_{\pi(k+1)}, \dots, x_{\pi(m)}) \quad (\pi \in \mathfrak{S}_m).$$

Then:

$$\begin{aligned} f \wedge g(x_{\pi(1)}, \dots, x_{\pi(m)}) &= \sum_{\pi' \in \mathfrak{S}_m} \text{sgn}(\pi') a(\pi' \circ \pi) \\ &= \text{sgn}(\pi) \sum_{\pi' \in \mathfrak{S}_m} \text{sgn}(\pi' \circ \pi) a(\pi' \circ \pi) \\ &= \text{sgn}(\pi) \sum_{\sigma \in \mathfrak{S}_m} \text{sgn}(\sigma) a(\sigma) \end{aligned}$$

Lemma 6.6 (associativity of the exterior product). *Let $f \in \wedge^k E$, $g \in \wedge^l E$ and $h \in \wedge^m E$ with $k, l, m \in \mathbb{N}_0$. Then we have:*

$$(f \wedge g) \wedge h = f \wedge (g \wedge h).$$

□

PROOF. Let $n := k + l + m$ and $\mathfrak{T} := \{\tau \in \mathfrak{S}_n : \tau(i) = i \text{ for } i > k + l\}$; moreover let for arbitrary $x_1, \dots, x_m \in E^*$ and $\pi \in \mathfrak{S}_m$

$$a(\pi) := f(x_{\pi(1)}, \dots, x_{\pi(k)}) g(x_{\pi(k+1)}, \dots, x_{\pi(k+l)}) h(x_{\pi(k+l+1)}, \dots, x_{\pi(n)}).$$

Herewith by twice applying Remark 6.5:

$$\begin{aligned} ((f \wedge g) \wedge h)(x_1, \dots, x_n) &= \\ &= \frac{1}{(k+l)!m!} \sum_{\sigma \in \mathfrak{S}_n} \text{sgn}(\sigma) \frac{1}{k!l!} \sum_{\tau \in \mathfrak{T}} \text{sgn}(\tau) a(\sigma \circ \tau) \\ &= \frac{1}{k!l!m!} \frac{1}{(k+l)!} \sum_{\tau \in \mathfrak{T}} \sum_{\sigma \in \mathfrak{S}_n} \text{sgn}(\sigma \circ \tau) a(\sigma \circ \tau) \\ &= \frac{1}{k!l!m!} \frac{\text{card}(\mathfrak{T})}{(k+l)!} \sum_{\gamma \in \mathfrak{S}_n} \text{sgn}(\gamma) a(\gamma) \\ &= \frac{1}{k!l!m!} \sum_{\gamma \in \mathfrak{S}_n} \text{sgn}(\gamma) a(\gamma). \end{aligned}$$

Since we obtain this result also, if we (with identical steps) calculate $(f \wedge (g \wedge h))(x_1, \dots, x_n)$, the claim is proven. □

Hence it is clear that expressions such as

$$f_1 \wedge \dots \wedge f_m \quad \text{with } f_l \in \wedge^{k_l} E$$

are uniquely determined. Thus we have:

Lemma 6.7. *Let $f_l \in \wedge^{k_l} E$ for $l \in \{1, \dots, m\}$. Then we have with $n := k_1 + \dots + k_m$*

$$f_1 \wedge \dots \wedge f_m = \prod_{1 \leq l \leq m} \frac{1}{k_l!} \cdot \sum_{\pi \in \mathfrak{S}_n} \text{sgn}(\pi) f_\pi,$$

where f_π with $i_l := k_1 + \dots + k_{l-1}$ is defined as

$$f_\pi(x_1, \dots, x_n) := f_1(x_{\pi(1)}, \dots, x_{\pi(i_1)}) f_2(x_{\pi(i_1+1)}, \dots, x_{\pi(i_2)}) \cdots f_m(x_{\pi(i_{m-1}+1)}, \dots, x_{\pi(n)}).$$

PROOF. This follows by induction on m . The case $l = 2$ is just the definition of 6.5; the case $l = 3$ is shown in the proof of 6.6. \square

Lemma 6.8. *Let e_1, \dots, e_d be a basis of $E^{**} \cong E$ and b_1, \dots, b_d the dual basis of E^* . Then we have for all $f \in \wedge^k E$:*

$$f = \sum_{i_1 < \dots < i_k} a_{i_1 \dots i_k} e_{i_1} \wedge \dots \wedge e_{i_k} \iff a_{i_1 \dots i_k} = f(b_{i_1}, \dots, b_{i_k}) \text{ for all } i_1 < \dots < i_k.$$

PROOF. Note first that for $i_1 < \dots < i_k$ and $j_1 < \dots < j_k$ from Lemma 6.7:

$$\begin{aligned} e_{i_1} \wedge \dots \wedge e_{i_k}(b_{j_1}, \dots, b_{j_k}) &= \sum_{\pi \in \mathfrak{S}_k} \text{sgn}(\pi) e_{i_1}(b_{j_{\pi(1)}}) \cdots e_{i_k}(b_{j_{\pi(k)}}) \\ &= \sum_{\pi \in \mathfrak{S}_k} \text{sgn}(\pi) \delta_{i_1, j_{\pi(1)}} \cdots \delta_{i_k, j_{\pi(k)}} \\ &= \delta_{i_1, j_1} \cdots \delta_{i_k, j_k} \quad (\text{da } i_1 < \dots < i_k) \\ &= \begin{cases} 1, & \text{if } i_1 = j_1, \dots, i_k = j_k \\ 0, & \text{else.} \end{cases} \end{aligned}$$

” \Rightarrow ” follows directly from this remark;

” \Leftarrow ” Let $g := \sum_{i_1 < \dots < i_k} f(b_{i_1}, \dots, b_{i_k}) e_{i_1} \wedge \dots \wedge e_{i_k} \in \wedge^k E$. By above remark we have $f(b_{i_1}, \dots, b_{i_k}) = g(b_{i_1}, \dots, b_{i_k})$ for all $i_1 < \dots < i_k$; by linearity therefore f and g are equal on $(E^*)^k$. \square

Lemma 6.9. *Let e_1, \dots, e_d be a basis of $E^{**} \cong E$ and $k \in \{1, \dots, d\}$. Then*

$$\{e_{i_1} \wedge \dots \wedge e_{i_k} : 1 \leq i_1 < \dots < i_k \leq d\}$$

is a basis of $\wedge^k E$. In particular we have:

$$\dim \wedge^k E = \binom{d}{k}.$$

PROOF. To prove this, we choose a basis b_1, \dots, b_d of E^* dual to e_1, \dots, e_d and apply 6.8. \square

Definition 6.10 (scalar product). Let b_1, \dots, b_d be a basis of E^* . Then

$$\langle f, g \rangle := \sum_{i_1 < \dots < i_k} f(b_{i_1}, \dots, b_{i_k}) g(b_{i_1}, \dots, b_{i_k}) \quad (f, g \in \wedge^k E)$$

defines a scalar product on $\wedge^k E$.

Lemma 6.11. For the scalar product from 6.10 we have:

$$\langle u_1 \wedge \dots \wedge u_k, v_1 \wedge \dots \wedge v_k \rangle = \det(\langle u_i, v_j \rangle)_{1 \leq i, j \leq k} \quad (u_i, v_i \in E).$$

In particular we have for the associated norm:

$$|u_1 \wedge \dots \wedge u_k| = \sqrt{\det(\langle u_i, u_j \rangle)_{1 \leq i, j \leq k}} \quad (u_i, v_i \in E).$$

PROOF. The right hand side $h(u_1, \dots, u_k; v_1, \dots, v_k) := \det(\langle u_i, v_j \rangle)_{1 \leq i, j \leq k}$ for fixed v_1, \dots, v_k is an alternating multilinear form in u_1, \dots, u_k and vice versa, i.e. $h(\cdot; v_1, \dots, v_k) \in \wedge^k E^*$ and $h(u_1, \dots, u_k; \cdot) \in \wedge^k E^*$. If e_1, \dots, e_d denotes the dual basis of b_1, \dots, b_d of E , and if we apply Lemma 6.8 twice, we obtain:

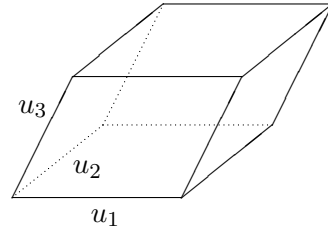
$$\begin{aligned} & h(u_1, \dots, u_k; v_1, \dots, v_k) \\ &= \sum_{i_1 < \dots < i_k} h(e_{i_1}, \dots, e_{i_k}; v_1, \dots, v_k) b_{i_1} \wedge \dots \wedge b_{i_k}(u_1, \dots, u_k) \\ &= \sum_{i_1 < \dots < i_k} \sum_{j_1 < \dots < j_k} h(e_{i_1}, \dots, e_{i_k}; e_{j_1}, \dots, e_{j_k}) b_{i_1} \wedge \dots \wedge b_{i_k}(u_1, \dots, u_k) \\ & \quad \cdot b_{j_1} \wedge \dots \wedge b_{j_k}(v_1, \dots, v_k) \\ &= \sum_{i_1 < \dots < i_k} b_{i_1} \wedge \dots \wedge b_{i_k}(u_1, \dots, u_k) b_{i_1} \wedge \dots \wedge b_{i_k}(v_1, \dots, v_k) \\ &= \sum_{i_1 < \dots < i_k} u_1 \wedge \dots \wedge u_k(b_{i_1}, \dots, b_{i_k}) v_1 \wedge \dots \wedge v_k(b_{i_1}, \dots, b_{i_k}) \\ &\equiv \langle u_1 \wedge \dots \wedge u_k, v_1 \wedge \dots \wedge v_k \rangle. \end{aligned}$$

□

Remark: In particular the preceding statement holds for $E := \mathbb{R}^d = E^*$ (with the canonical basis). The norm

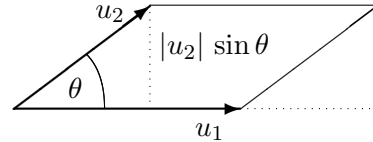
$$|u_1 \wedge \dots \wedge u_k| = \sqrt{\det(\langle u_i, u_j \rangle)_{1 \leq i, j \leq k}}$$

provides for $E = \mathbb{R}^d$ the k -dimensional volume of the parallelepiped spanned by u_1, \dots, u_k .



So we have e.g. in the case $k = 2$:

$$\begin{aligned}
 |u_1 \wedge u_2| &= \det \begin{pmatrix} \langle u_1, u_1 \rangle & \langle u_1, u_2 \rangle \\ \langle u_2, u_1 \rangle & \langle u_2, u_2 \rangle \end{pmatrix}^{1/2} \\
 &= \det \begin{pmatrix} |u_1|^2 & |u_1| |u_2| \cos \theta \\ |u_1| |u_2| \cos \theta & |u_2|^2 \end{pmatrix}^{1/2} \\
 &= (|u_1|^2 |u_2|^2 (1 - \cos^2 \theta))^{1/2} \\
 &= |u_1| |u_2| |\sin \theta|
 \end{aligned}$$



□

Since we want to study the action of A on k -dimensional objects in \mathbb{R}^d , we now consider the k -fold exterior product of a matrix:

Definition-remark 6.12. Let $A \in \mathbb{R}^{d \times d}$. Then by Lemma 6.9 via

$$\wedge^k A (u_1 \wedge \dots \wedge u_k) := Au_1 \wedge \dots \wedge Au_k \quad (u_i \in \mathbb{R}^d)$$

a linear operator $\wedge^k A : \wedge^k \mathbb{R}^d \rightarrow \wedge^k \mathbb{R}^d$ is defined, the k -fold exterior product of the matrix A . For this we have:

- i) $\wedge^1 A = A$,
- ii) $\wedge^d A = \det A$ (by Lemma 6.7),
- iii) $\wedge^k(AB) = (\wedge^k A)(\wedge^k B)$,
- iv) $(\wedge^k A)^{-1} = \wedge^k A^{-1}$ if A invertible,
- v) $\wedge^k(cA) = c^k \wedge^k A$ for $c \in \mathbb{R}$,
- vi) $\wedge^k U$ orthogonal, if U orthogonal and in this case we have $(\wedge^k U)^* = \wedge^k U^*$.

Lemma 6.13 (exterior product of a matrix and eigenvalues). *Let $\lambda_1, \dots, \lambda_d$ be the eigenvalues of $A \in \mathbb{R}^{d \times d}$. Then $\wedge^k A$ has the eigenvalues*

$$\{ \lambda_{i_1} \cdots \lambda_{i_k} : 1 \leq i_1 < \dots < i_k \leq d \}.$$

PROOF. If u_1, \dots, u_d are the eigenvectors of $\lambda_1, \dots, \lambda_d$, fixing indices $1 \leq i_1 < \dots < i_k \leq d$, we get:

$$\begin{aligned} \wedge^k A (u_{i_1} \wedge \dots \wedge u_{i_k}) &\equiv Au_{i_1} \wedge \dots \wedge Au_{i_k} \\ &= \lambda_{i_1} u_{i_1} \wedge \dots \wedge \lambda_{i_k} u_{i_k} \\ &= (\lambda_{i_1} \cdots \lambda_{i_k}) (u_{i_1} \wedge \dots \wedge u_{i_k}), \end{aligned}$$

so that $\lambda_{i_1} \cdots \lambda_{i_k}$ is an eigenvalue with eigenvector $u_{i_1} \wedge \dots \wedge u_{i_k}$. For dimension reasons these are all eigenvectors and hence all eigenvalues. \square

Lemma 6.14 (exterior product of a matrix and singular value decomposition). *For $A \in \mathbb{R}^{d \times d}$ let $\delta_1 \geq \dots \geq \delta_d \geq 0$ be the singular values and*

$$A = VDU$$

a singular value decomposition, where $D \equiv \text{diag}(\delta_1, \dots, \delta_d)$. Then we have for $k = 1, \dots, d$:

- i) $\wedge^k A = (\wedge^k V)(\wedge^k D)(\wedge^k U)$ is singular value decomposition of $\wedge^k A$;
- ii) $\wedge^k D = \text{diag}(\delta_{i_1} \cdots \delta_{i_k} : 1 \leq i_1 < \dots < i_k \leq d)$.
Hence $\delta_1 \cdots \delta_k$ is the biggest resp. $\delta_{d-k+1} \cdots \delta_d$ the smallest singular value of $\wedge^k A$.
- iii) For the operator norm we have:
 $\| \wedge^k A \| = \delta_1 \cdots \delta_k$, $|\det A| = \| \wedge^d A \| = \delta_1 \cdots \delta_d$ and $\| \wedge^k A \| \leq \| A \|^k$.

PROOF. i) and ii) follow from Remark 6.12 and Lemma 6.13; iii) follows from ii) and the definition of the operator norm $\| \cdot \|$. \square

Theorem 6.15 (Furstenberg-Kesten). *Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space and $A : \Omega \rightarrow \mathbb{R}^{d \times d}$ a random matrix, for which we have*

$$\log^+ \|A(\cdot)\| \in L^1(\mathbb{P}). \quad (3)$$

Moreover let as in (2)

$$A_n := (A \circ \varphi^{n-1}) (A \circ \varphi^{n-2}) \cdots (A \circ \varphi) A$$

with a (\mathbb{P}) -measure preserving map $\varphi : \Omega \rightarrow \Omega$.

Then there exists a set $\tilde{\Omega} \in \mathcal{F}$ with $\mathbb{P}(\tilde{\Omega}) = 1$ and $\tilde{\Omega} \subset \varphi^{-1}(\tilde{\Omega})$, and there exist measurable functions

$$\gamma^{(k)} : \Omega \longrightarrow \mathbb{R} \cup \{-\infty\} \quad (k = 1, \dots, d)$$

with $\gamma^{(k)+} \in L^1(\mathbb{P})$, such that for all $\omega \in \tilde{\Omega}$ and $k, m \in \{1, \dots, d\}$ we have:

$$\begin{aligned} \gamma^{(k)}(\omega) &= \lim_{n \rightarrow \infty} \frac{1}{n} \log \|\wedge^k A_n(\omega)\|, \\ \gamma^{(k)}(\varphi(\omega)) &= \gamma^{(k)}(\omega), \\ \gamma^{(k+m)}(\omega) &\leq \gamma^{(k)}(\omega) + \gamma^{(m)}(\omega). \end{aligned}$$

If we define recursively random variables

$$\Lambda_k : \Omega \longrightarrow \mathbb{R} \cup \{-\infty\} \quad (k = 1, \dots, d)$$

by

$$\Lambda_1 + \dots + \Lambda_k = \gamma^{(k)}$$

with

$$\Lambda_k := -\infty \quad \text{on} \quad \{\gamma^{(k)} = -\infty\},$$

we have for all $\omega \in \tilde{\Omega}$ and $k \in \{1, \dots, d\}$:

$$\begin{aligned} \Lambda_k(\omega) &= \lim_{n \rightarrow \infty} \frac{1}{n} \log \delta_k(A_n(\omega)), \\ \Lambda_k(\varphi(\omega)) &= \Lambda_k(\omega), \\ \Lambda_1(\omega) &\geq \Lambda_2(\omega) \geq \dots \geq \Lambda_d(\omega) \quad (\geq -\infty). \end{aligned}$$

If \mathbb{P} is ergodic, so $\gamma^{(k)}$ and Λ_k are constant by the invariance above (on $\tilde{\Omega}$), hence $\gamma^{(k)} = \mathbb{E}(\gamma^{(k)})$ and $\Lambda_k = \mathbb{E}(\Lambda_k)$.

PROOF. 1) Let

$$Y_n^k := \log \|\wedge^k A_n\| \quad (n \in \mathbb{N}, k = 1, \dots, d);$$

then $(Y_n^k)_n$ for each $k = 1, \dots, d$ is subadditive: in case $k = 1$ this had been shown in 5.3; for $k > 1$ the calculation can immediately be transferred, since for all matrices B, C we have: $\wedge^k(BC) = (\wedge^k B)(\wedge^k C)$. Consequently with A also each $\wedge^k A$ is a cocycle, i.e. we have:

$$\wedge^k A_{n+m} = \wedge^k A_n \circ \varphi^m \cdot \wedge^k A_m.$$

Hence subadditivity of $(Y_n^k)_n$ follows.

2) The existence of $\tilde{\Omega}$ and $\gamma^{(k)}$ with the claimed properties follows from Theorem 5.7, applied to $(-Y_n^k)_n$; it remains to prove:

$$\gamma^{(k+m)} \leq \gamma^{(k)} + \gamma^{(m)};$$

but this follows directly from the characteristic property of the $\gamma^{(k)}$ and the norm inequality

$$\|\wedge^{k+m} A_n\| \leq \|\wedge^k A_n\| \cdot \|\wedge^m A_n\|.$$

3) We now prove the claims with respect to Λ_k : By Lemma 6.14 we have for $k = 1, \dots, d$:

$$\frac{1}{n} \log \|\wedge^k A_n\| = \frac{1}{n} \sum_{i=1}^k \log \delta_i(A_n).$$

We have $\Lambda_1 \equiv \gamma^{(1)}$ and for $\omega \in \tilde{\Omega}$ we obtain successively:

$$\Lambda_{k+1}(\omega) \equiv \gamma^{(k+1)}(\omega) - \gamma^{(k)}(\omega) = \lim_{n \rightarrow \infty} \frac{1}{n} \log \delta_{k+1}(A_n),$$

if $\gamma^k(\omega) > -\infty$; if this procedure ends, i.e. if $\gamma^{k_0}(\omega) = -\infty$, then also $\gamma^k(\omega) = -\infty$ for all $k \geq k_0$ and thus also $\Lambda_k = -\infty$ for all $k \geq k_0$. The remaining statements hold true by

$$\delta_1(A_n) \geq \delta_2(A_n) \geq \dots \geq \delta_d(A_n)$$

and the respective expectations exist by hypothesis. □

7. The multiplicative ergodic theorem of Oseledets

Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space with a measure preserving mapping $\varphi : \Omega \rightarrow \Omega$ and $A : \Omega \rightarrow \mathbb{R}^{d \times d}$ a random matrix. According to (2) we define

$$A_n := \begin{cases} (A \circ \varphi^{n-1}) (A \circ \varphi^{n-2}) \cdots (A \circ \varphi) A, & n \in \mathbb{N}, \\ I, & n = 0, \end{cases}$$

the *cocycle generated by A*; A_n is therefore a *cocycle* over φ , i.e. we have:

$$A_{n+m} = (A_n \circ \varphi^m) \cdot A_m \quad (m, n \in \mathbb{N}_0),$$

as has already been used in the proof of the theorem of Furstenberg-Kesten.

We are now interested in the asymptotics of $|A_n x|$ for $x \in \mathbb{R}^d$ as $n \rightarrow \infty$. We trace this behavior back to the theorem of Furstenberg-Kesten by means of the following (deterministic) proposition 7.3. To prove the convergence statements therein we first establish two Lemmata:

Lemma 7.1. *Let $\Phi \in \mathbb{R}^{d \times d}$ be symmetric with spectral decomposition*

$$\Phi = \sum_{i=1}^r \lambda_i P_i,$$

where $r \leq d$ and λ_i the eigenvalues and P_i the corresponding orthogonal projectors on the eigenspaces. Let

$$\Phi_n = \sum_{i=1}^{r_n} \lambda_i^n P_i^n$$

equally symmetric $d \times d$ -matrices, so that we have:

- i) $\lambda_k^n \xrightarrow{n \rightarrow \infty} \lambda_i$ for all $k \in \Sigma_i$, where $\Sigma_i \neq \emptyset$ are sets of indices ($i=1, \dots, r$);
- ii) $\bar{P}_i^n := \sum_{k \in \Sigma_i} P_k^n \xrightarrow{n \rightarrow \infty} P_i$ for all $i = 1, \dots, r$.

Then: $\Phi_n \xrightarrow{n \rightarrow \infty} \Phi$.

PROOF. With the convergence statements we obtain:

$$\begin{aligned} \Phi_n - \Phi &= \sum_{i=1}^r \sum_{k \in \Sigma_i} \lambda_k^n P_k^n - \sum_{i=1}^r \lambda_i P_i \\ &= \sum_{i=1}^r \left[\underbrace{\sum_{k \in \Sigma_i} (\lambda_k^n - \lambda_i) P_k^n}_{\rightarrow 0} + \lambda_i \underbrace{\left(\sum_{k \in \Sigma_i} P_k^n - P_i \right)}_{\rightarrow 0} \right] \xrightarrow{n \rightarrow \infty} 0. \end{aligned}$$

□

Lemma 7.2. *Let P, Q be orthogonal projectors in \mathbb{R}^2 , such that we have:*

$$\dim U = \dim V = 1, \quad \text{where } U := \text{Im } P \text{ and } V := \text{Im } Q.$$

Then:

$$\delta(U, V) := \|P - Q\| = |x \wedge y| = |\sin \alpha| \quad (x \in U, y \in V \text{ with } |x| = |y| = 1),$$

where α is the angle between x and y . Consequently δ is a complete metric on \mathbb{P}^1 , the projective space of all one dimensional linear subspaces of \mathbb{R}^2 .

PROOF. The second equation has already been proved in the remark after Lemma 6.11.

$\|P - Q\| = |x \wedge y|$: As in the remark after Lemma 6.11 we further obtain:

$$\begin{aligned} |x \wedge y| &= \det \begin{pmatrix} \langle x, x \rangle & \langle x, y \rangle \\ \langle y, x \rangle & \langle y, y \rangle \end{pmatrix}^{1/2} \\ &= \sqrt{1 - \langle x, y \rangle^2} \\ &= \sqrt{\langle x, y \rangle^2 + \langle x, y^\perp \rangle^2 - \langle x, y \rangle^2} \\ &= |\langle x, y^\perp \rangle| \\ &= \|(I - Q)P\| \\ &= \|(P - Q)P\| \leq \|P - Q\|, \end{aligned}$$

where the idempotence of orthogonal projectors has been used besides the fact $\|AB\| = \|BA\|$ for orthogonal projectors A, B .

Note that $\|PQ\| = |\langle x, y \rangle| = |\cos(\alpha)| = |\langle x^\perp, y^\perp \rangle| = \|(I - Q)(I - P)\|$.

On the other hand for $w \in \mathbb{R}^2$:

$$\begin{aligned} |(P - Q)w|^2 &= |(P - QP)w - (Q - QP)w|^2 \\ &= |(I - Q)Pw - Q(I - P)w|^2 \\ &= |(I - Q)Pw|^2 + |Q(I - P)w|^2 \\ &\leq \|(I - Q)P\|^2 |Pw|^2 + \underbrace{\|Q(I - P)\|^2}_{\|(I - Q)P\|^2} |(I - P)w|^2 \\ &= \|(I - Q)P\|^2, \end{aligned}$$

hence

$$\|P - Q\| \leq \|(I - Q)P\|.$$

In summary we proved:

$$\|P - Q\| = \|(I - Q)P\| = |x \wedge y|.$$

□

The following deterministic theorem serves to prepare for an application of the theorem of Furstenberg-Kesten.

Proposition 7.3 (Goldsheid-Margulis). *Let $(A_n)_{n \in \mathbb{N}}$ be a sequence in $\mathbb{R}^{d \times d}$ with the properties:*

$$\limsup_{n \rightarrow \infty} \frac{1}{n} \log \|A_n\| \leq 0 \quad (4)$$

and assume that $\Phi_n := A_n \cdots A_1$ fulfils

$$\lim_{n \rightarrow \infty} \frac{1}{n} \log \|\wedge^i \Phi_n\| =: \gamma^{(i)} \in \mathbb{R} \cup \{-\infty\} \quad (5)$$

for each $i = 1, \dots, d$. Then we have:

i) There exists (in the topology of the operator norm) the limit

$$\Psi := \lim_{n \rightarrow \infty} (\Phi_n^* \Phi_n)^{1/2n} \geq 0.$$

Defining successively Λ_i for $i = 1, \dots, d$ by $\Lambda_1 + \dots + \Lambda_i = \gamma^{(i)}$ (if $\gamma^{(i)} = -\infty$, set $\Lambda_i = -\infty$), then the eigenvalues of Ψ are given by

$$e^{\Lambda_1}, \dots, e^{\Lambda_d}$$

and we have

$$\Lambda_i = \lim_{n \rightarrow \infty} \frac{1}{n} \log \delta_i(\Phi_n) \quad (i = 1, \dots, d).$$

ii) Let

$$e^{\lambda_p} < \dots < e^{\lambda_1}$$

the different (!) eigenvalues of Ψ (where $\lambda_p = -\infty$ is possible), U_p, \dots, U_1 the corresponding eigenspaces with $d_i := \dim U_i$ and let

$$V_i := \begin{cases} \{0\}, & i = p+1 \\ U_p \oplus \dots \oplus U_i, & i = 1, \dots, p. \end{cases}$$

Then we have:

$$V_{p+1} \subset V_p \subset V_{p-1} \subset \dots \subset V_1 = \mathbb{R}^d$$

and for each $x \in \mathbb{R}^d \setminus \{0\}$ there exists the Lyapunov exponent

$$\lambda(x) := \lim_{n \rightarrow \infty} \frac{1}{n} \log |\Phi_n x|;$$

we have for all $i = 1, \dots, p$:

$$x \in V_i \setminus V_{i+1} \iff \lambda(x) = \lambda_i$$

resp. equivalently:

$$V_i = \{x \in \mathbb{R}^d : \lambda(x) \leq \lambda_i\}.$$

PROOF. In case $d = 1$ nothing needs to be proven, since then $\Phi_n \in \mathbb{R}$ and the claims follow directly from the hypotheses.

For simplicity we now confine our attention to the case $d = 2$; the general case can be proved in a similar way, with more technicalities (see Arnold [AR 98] pp. 144-152).

$\Lambda_i = \lim_{n \rightarrow \infty} \frac{1}{n} \log \delta_i(\Phi_n)$ for $i = 1, 2$ This follows from (5) with Lemma 6.14 iii):

$$\Lambda_1 \equiv \gamma^1 \stackrel{(5)}{=} \lim_{n \rightarrow \infty} \frac{1}{n} \log \|\Phi_n\| \stackrel{6.14}{=} \lim_{n \rightarrow \infty} \frac{1}{n} \log \delta_1(\Phi_n);$$

if now $\Lambda_1 = -\infty$, hence $\gamma^1 = -\infty$, by (5) also $\gamma^2 = -\infty = \Lambda_2$; on the other hand in this case

$$\frac{1}{n} \log \delta_2(\Phi_n) \leq \frac{1}{n} \log \delta_1(\Phi_n) \longrightarrow -\infty.$$

If $\Lambda_1 > -\infty$, we get:

$$\begin{aligned} \Lambda_2 \equiv \gamma^2 - \Lambda_1 &= \lim_{n \rightarrow \infty} \frac{1}{n} \log \underbrace{\|\wedge^2 \Phi_n\|}_{\delta_1(\Phi_n)\delta_2(\Phi_n)} - \lim_{n \rightarrow \infty} \frac{1}{n} \log \delta_1(\Phi_n) \\ &= \lim_{n \rightarrow \infty} \frac{1}{n} \log \delta_2(\Phi_n). \end{aligned}$$

convergence of operators and Lyapunov exponents Let now

$$\Phi_n = V_n D_n O_n$$

be the singular value decomposition of Φ_n , with

$$D_n = \begin{pmatrix} \delta_1(\Phi_n) & 0 \\ 0 & \delta_2(\Phi_n) \end{pmatrix}.$$

from this we obtain:

$$(\Phi_n^* \Phi_n)^{1/2n} = (O_n^* D_n^2 O_n)^{1/2n} = O_n^* D_n^{1/n} O_n;$$

this matrix has eigenvalues $\delta_1(\Phi_n)^{1/n}$ and $\delta_2(\Phi_n)^{1/n}$, which according to what has been shown above converge to e^{Λ_1} and e^{Λ_2} ; so we have the following convergences:

$$D_n^{1/n} \equiv \begin{pmatrix} \delta_1^{1/n}(\Phi_n) & 0 \\ 0 & \delta_2^{1/n}(\Phi_n) \end{pmatrix} \xrightarrow{n \rightarrow \infty} \begin{pmatrix} e^{\Lambda_1} & 0 \\ 0 & e^{\Lambda_2} \end{pmatrix}.$$

Now the difficulty is that the convergence of O_n in general is not guaranteed; but it is enough to prove convergence of the respective eigenspaces for which Lemma 7.1 has been established.

1. CASE: $\Lambda_1 = \Lambda_2 =: \lambda_1$: As just seen we have $D_n^{1/n} \rightarrow e^{\lambda_1} I$ and Lemma 7.1 imply

$$(\Phi_n^* \Phi_n)^{1/2n} \xrightarrow{n \rightarrow \infty} e^{\lambda_1} I$$

with $\bar{P}_1^n := P_1^n + P_2^n$. Moreover we have immediately: $V_1 \equiv U_1 = \mathbb{R}^2$, $p = 1$ and $d_1 = 2$. Therefore we only have to prove that for all $x \in \mathbb{R}^2 \setminus \{0\}$ we have:

$$\lambda(x) \equiv \lim_{n \rightarrow \infty} \frac{1}{n} \log \|\Phi_n x\| = \lambda_1.$$

For this let first $\lambda_1 > -\infty$; then from the already proven characterization of Λ_1 , it follows that for each $\epsilon > 0$ there exists $c_\epsilon \in (0, \infty)$ with

$$\frac{1}{c_\epsilon} e^{n(\lambda_1 - \epsilon)} \leq \delta_i(\Phi_n) \leq c_\epsilon e^{n(\lambda_1 + \epsilon)}, \quad i = 1, 2.$$

Setting $x_n := O_n x$, we get

$$|\Phi_n x| = |V_n D_n O_n x| = |D_n x_n| = (\delta_1(\Phi_n)^2 (x_n^1)^2 + \delta_2(\Phi_n)^2 (x_n^2)^2)^{1/2}$$

with x_n^i denoting the components of x_n ; in summary we therefore have

$$\frac{|x|}{c_\epsilon} e^{n(\lambda_1 - \epsilon)} \leq |\Phi_n x| \leq |x| c_\epsilon e^{n(\lambda_1 + \epsilon)},$$

whence we obtain that $\lambda(x) = \lambda_1$.

If $\lambda = -\infty$, we can in the same way find for each $r < 0$ a $c_r \in (0, \infty)$ such that

$$0 \leq \delta_i(\Phi_n) \leq c_r e^{nr}, \quad i = 1, 2.$$

As above we then infer:

$$0 \leq |\Phi_n x| \leq |x| c_r e^{nr},$$

from which we conclude as above: $\lambda(x) = \lambda_1$. So the theorem is proved in case $\Lambda_1 = \Lambda_2$.

2. CASE: $\lambda_1 \equiv \Lambda_1 > \Lambda_2 \equiv \lambda_2$: Here we have

$$D_n^{1/n} \equiv \begin{pmatrix} \delta_1^{1/n}(\Phi_n) & 0 \\ 0 & \delta_2^{1/n}(\Phi_n) \end{pmatrix} \xrightarrow{n \rightarrow \infty} \begin{pmatrix} e^{\lambda_1} & 0 \\ 0 & e^{\lambda_2} \end{pmatrix}.$$

To prove the existence of Ψ , we have to show that the orthogonal projectors P_1^n, P_2^n on the eigenspaces U_1^n, U_2^n of $(\Phi_n^* \Phi_n)^{1/2n}$ converge to orthogonal projectors P_1, P_2 , since then Lemma 7.1 yields:

$$(\Phi_n^* \Phi_n)^{1/2n} \xrightarrow{n \rightarrow \infty} e^{\lambda_1} P_1 + e^{\lambda_2} P_2 =: \Psi.$$

This will be proved in the following Lemma by means of a Cauchy sequence argument. For this purpose we remark that the eigenvectors of $(\Phi_n^* \Phi_n)^{1/2n} = O_n^* D_n^{1/n} O_n$ are given by $u_i^n := O_n^* e_i$ ($i = 1, 2$), where (e_1, e_2) is the standard basis of \mathbb{R}^2 . In particular $U_i^n = \text{span}(u_i^n)$, $i = 1, 2$.

Lemma 7.4. *In the situation above ("2. case" in the proof of theorem 7.3) we have:*

$$\limsup_{n \rightarrow \infty} \frac{1}{n} \log \delta(U_i^n, U_i^{n+1}) \leq \lambda_2 - \lambda_1 < 0 \quad (i = 1, 2).$$

In particular $(U_i^n)_{n \in \mathbb{N}}$ ($i = 1, 2$) is a Cauchy sequence in the projective space P^1 , that hence converges to $U_i \in P^1$. Moreover, this convergence takes place with exponential speed:

$$\limsup_{n \rightarrow \infty} \frac{1}{n} \log \delta(U_i^n, U_i) \leq \lambda_2 - \lambda_1 \quad (i = 1, 2).$$

PROOF OF LEMMA 7.4. Wlog let hereby $i = 2$, since U_1^n is orthogonal to U_2^n , but the metric δ on \mathbb{P}^1 is invariant with respect to orthogonal transformations. By orthogonality of all (u_1^{n+1}, u_2^{n+1}) we may represent u_2^n as

$$u_2^n = \alpha_n u_1^{n+1} + \beta_n u_2^{n+1} \quad (n \in \mathbb{N}).$$

1) $\delta(U_2^n, U_2^{n+1}) = |\alpha_n|$, because:

$$\begin{aligned} \delta(U_2^n, U_2^{n+1}) &\stackrel{7.2}{=} |u_2^n \wedge u_2^{n+1}| \equiv |(\alpha_n u_1^{n+1} + \beta_n u_2^{n+1}) \wedge u_2^{n+1}| \\ &= |\alpha_n| |u_1^{n+1} \wedge u_2^{n+1}| \\ &= |\alpha_n|, \end{aligned}$$

where orthonormality of u_1^n and u_2^n was used.

2) $\delta(U_2^n, U_2^{n+1}) \leq \|A_{n+1}\| \frac{\delta_2(\Phi_n)}{\delta_1(\Phi_{n+1})}$, because: First we have

$$\begin{aligned} |\Phi_{n+1} u_2^n| &\equiv |\alpha_n \Phi_{n+1} u_1^{n+1} + \beta_n \Phi_{n+1} u_2^{n+1}| \\ &\equiv |\alpha_n V_{n+1} D_{n+1} O_{n+1} O_{n+1}^* e_1 + \beta_n V_{n+1} D_{n+1} O_{n+1} O_{n+1}^* e_2| \\ &= |\alpha_n \delta_1(\Phi_{n+1}) V_{n+1} e_1 + \beta_n \delta_2(\Phi_{n+1}) V_{n+1} e_2| \\ &\stackrel{\text{orth.}}{\geq} |\alpha_n \delta_1(\Phi_{n+1}) V_{n+1} e_1| \\ &= |\alpha_n| \delta_1(\Phi_{n+1}); \end{aligned}$$

on the other hand

$$|\Phi_{n+1} u_2^n| \equiv |A_{n+1} \Phi_n u_2^n| \leq \|A_{n+1}\| |\Phi_n u_2^n| = \|A_{n+1}\| \delta_2(\Phi_n),$$

hence in summary

$$\delta(U_2^n, U_2^{n+1}) \stackrel{1)}{=} |\alpha_n| \leq \frac{|\Phi_{n+1} u_2^n|}{\delta_1(\Phi_{n+1})} \leq \|A_{n+1}\| \frac{\delta_2(\Phi_n)}{\delta_1(\Phi_{n+1})}.$$

3) First claim of the Lemma: by what has just been proved we get:

$$\begin{aligned} \limsup_{n \rightarrow \infty} \frac{1}{n} \log \delta(U_2^n, U_2^{n+1}) &\leq \limsup_{n \rightarrow \infty} \frac{1}{n} \log \|A_{n+1}\| \\ &\quad + \limsup_{n \rightarrow \infty} \frac{1}{n} \log \delta_2(\Phi_n) \\ &\quad - \liminf_{n \rightarrow \infty} \frac{1}{n} \log \delta_1(\Phi_{n+1}) \\ &\leq 0 + \lambda_2 - \lambda_1, \end{aligned}$$

where we use the first hypothesis of theorem 7.3 and the convergence result already established.

4) $(U_2^n)_n$ converges in P^1 to some U_2 : For this purpose let $\varepsilon < \lambda_1 - \lambda_2$; by what has been shown we can choose $n_0 \in \mathbb{N}$ such that

$$\frac{1}{n} \log \delta(U_2^n, U_2^{n+1}) < \lambda_2 - \lambda_1 + \varepsilon \quad (< 0) \quad (\forall n \geq n_0).$$

But then we get for $n_0 \leq m \leq n$:

$$\begin{aligned} \delta(U_2^n, U_2^{n+1}) &\leq \sum_{k=m}^{n-1} \delta(U_2^k, U_2^{k+1}) \\ &\leq \sum_{k=m}^{n-1} e^{k(\lambda_2 - \lambda_1 + \varepsilon)} \\ &\leq \sum_{k=m}^{\infty} e^{k(\lambda_2 - \lambda_1 + \varepsilon)} \\ &= \frac{e^{m(\lambda_2 - \lambda_1 + \varepsilon)}}{1 - e^{\lambda_2 - \lambda_1 + \varepsilon}} \xrightarrow{m \rightarrow \infty} 0, \end{aligned}$$

where the summation formula for geometric series was employed.

5) Second claim of the Lemma: with the arguments just used we also get:

$$\delta(U_2^n, U_2) \leq e^{n(\lambda_2 - \lambda_1 + \varepsilon)} \frac{1}{1 - e^{\lambda_2 - \lambda_1 + \varepsilon}}$$

and therefore

$$\limsup_{n \rightarrow \infty} \frac{1}{n} \log \delta(U_2^n, U_2) \leq \lambda_2 - \lambda_1 + \varepsilon.$$

Now the claim follows with $\varepsilon \rightarrow 0$.

□
Lem.
7.4

CONTINUATION OF THE PROOF OF PROPOSITION 7.3. As orthogonal projectors P_1, P_2 we now choose the projectors on the spaces U_1, U_2 existing due to Lemma 7.4. By Lemma 7.2 and Lemma 7.4 we get the convergence

$$P_i^n \xrightarrow{n \rightarrow \infty} P_i \quad (i = 1, 2).$$

In summary we obtain

$$(\Phi_n^* \Phi_n)^{1/2n} \xrightarrow{n \rightarrow \infty} e^{\lambda_1} P_1 + e^{\lambda_2} P_2 =: \Psi$$

It remains to prove the claim on the *Lyapunov exponents*; hereby $V_2 = U_2 \subset \mathbb{R}^2 = V_1$, such that it remains to prove:

$$\begin{aligned} x \in V_2 \setminus \{0\} &\implies \lim_{n \rightarrow \infty} \frac{1}{n} \log \|\Phi_n x\| = \lambda_2 \quad \text{and} \\ x \in \mathbb{R}^2 \setminus V_2 &\implies \lim_{n \rightarrow \infty} \frac{1}{n} \log \|\Phi_n x\| = \lambda_1; \end{aligned}$$

where in each case we may assume $|x| = 1$.

$x \in V_2 \setminus \{0\} \implies \lim_{n \rightarrow \infty} \frac{1}{n} \log |\Phi_n x| = \lambda_2$: We represent x as

$$x = \alpha_n u_1^n + \beta_n u_2^n,$$

hence again

$$\Phi_n x = \alpha_n \Phi_n u_1^n + \beta_n \Phi_n u_2^n = \alpha_n \delta_1(\Phi_n) V_n e_1 + \beta_n \delta_2(\Phi_n) V_n e_2,$$

and thus

$$|\beta_n| \delta_2(\Phi_n) \leq [\alpha_n^2 \delta_1(\Phi_n)^2 + \beta_n^2 \delta_2(\Phi_n)^2]^{1/2} = |\Phi_n x|;$$

as in the proof of 1) of Lemma 7.4 we obtain from Lemma 7.2: $\delta(U_2^n, U_2^{n+1}) = |\alpha_n|$, since $x \in V_2 = U_2$; consequently by Lemma 7.4 also

$$\limsup_{n \rightarrow \infty} \frac{1}{n} \log |\alpha_n| = \limsup_{n \rightarrow \infty} \frac{1}{n} \log \delta(U_2^n, U_2) \leq \lambda_2 - \lambda_1 < 0;$$

thus we infer:

$$\beta_n^2 = 1 - \alpha_n^2 \xrightarrow{n \rightarrow \infty} 1.$$

and therefore in summary:

$$\begin{aligned} \lambda_2 &= \lim_{n \rightarrow \infty} \frac{1}{n} \log (|\beta_n| \delta_2(\Phi_n)) \\ &\leq \liminf_{n \rightarrow \infty} \frac{1}{n} \log |\Phi_n x| \leq \limsup_{n \rightarrow \infty} \frac{1}{n} \log |\Phi_n x| \\ &= \frac{1}{2} \limsup_{n \rightarrow \infty} \frac{1}{n} \log [\alpha_n^2 \delta_1(\Phi_n)^2 + \beta_n^2 \delta_2(\Phi_n)^2] \\ &\leq \frac{1}{2} \max \left\{ \limsup_{n \rightarrow \infty} \frac{1}{n} \log \alpha_n^2 \delta_1(\Phi_n)^2, \limsup_{n \rightarrow \infty} \frac{1}{n} \log \beta_n^2 \delta_2(\Phi_n)^2 \right\} \\ &\leq \max \{ (\lambda_2 - \lambda_1) + \lambda_1, 0 + \lambda_2 \} \\ &= \lambda_2. \end{aligned}$$

$x \in \mathbb{R}^2 \setminus V_2 \Rightarrow \lim_{n \rightarrow \infty} \frac{1}{n} \log |\Phi_n x| = \lambda_1$: Here we represent x as

$$x = \alpha u + \beta v$$

with unit vectors $u \in U_1$ and $v \in U_2 = V_2$; the latter are represented by

$$v = \alpha_n u_1^n + \beta_n u_2^n \quad \text{resp.} \quad u = \gamma_n u_1^n + \delta_n u_2^n .$$

Also in this case we necessarily get from Lemma 7.4: $\alpha_n \rightarrow 0$, $\delta_n \rightarrow 0$ and thus $|\beta_n| \rightarrow 1$, $|\gamma_n| \rightarrow 1$ (in projective space we have by Lemma 7.4: $u_1^n \rightarrow u$ and $u_2^n \rightarrow v$).

Therefore we have as above:

$$\begin{aligned} |\alpha| |\gamma_n| \delta_1(\Phi_n) &\leq [(\alpha \gamma_n + \beta \alpha_n)^2 \delta_1(\Phi_n)^2 + (\alpha \delta_n + \beta \beta_n)^2 \delta_2(\Phi_n)^2]^{1/2} \\ &= |\Phi_n x| ; \end{aligned}$$

noting that by the position of x always $\alpha = \langle x, u \rangle \neq 0$, we in summary again obtain:

$$\begin{aligned} \lambda_1 &= \lim_{n \rightarrow \infty} \frac{1}{n} \log (|\alpha| |\gamma_n| \delta_1(\Phi_n)) \\ &\leq \liminf_{n \rightarrow \infty} \frac{1}{n} \log |\Phi_n x| \leq \limsup_{n \rightarrow \infty} \frac{1}{n} \log |\Phi_n x| \\ &= \frac{1}{2} \limsup_{n \rightarrow \infty} \frac{1}{n} \log [(\alpha \gamma_n + \beta \alpha_n)^2 \delta_1(\Phi_n)^2 + (\alpha \delta_n + \beta \beta_n)^2 \delta_2(\Phi_n)^2] \\ &\leq \lambda_1 . \end{aligned}$$

Thus all claims of Proposition 7.3 have been proven. \square

To be able to apply the Proposition by Goldsheid-Margulis, it remains to check the first hypothesis in the special case of stationary random matrices:

Lemma 7.5. *Let $X : \Omega \rightarrow \mathbb{R} \cup \{-\infty\}$ be a random variable with $X^+ \in L^1(\Omega, \mathcal{F}, \mathbb{P})$. Then*

$$\Omega_1 := \left\{ \limsup_{n \rightarrow \infty} \frac{1}{n} X \circ \varphi^{n-1} \leq 0 \right\}$$

is invariant and we have $\mathbb{P}(\Omega_1) = 1$.

PROOF. The invariance follows from the definition of Ω_1 . Moreover Ω_1 has full measure, because:

$$\begin{aligned} \sum_{n=1}^{\infty} \mathbb{P} \left\{ \frac{1}{n} X \circ \varphi^{n-1} > \varepsilon \right\} &\stackrel{\varphi \text{ m. pres.}}{=} \sum_{n=1}^{\infty} \mathbb{P} \{X > \varepsilon n\} = \sum_{n=1}^{\infty} \mathbb{P} \{X^+ > \varepsilon n\} \\ &\leq \frac{1}{\varepsilon} \mathbb{E}(X^+) < \infty , \end{aligned}$$

hence by Borel-Cantelli: $\mathbb{P}(\Omega_1) = 1$. \square

To deduce the main theorem, we apply Lemma 7.5 to $X := \log \|A\|$. So we obtain:

Theorem 7.6 (Multiplicative ergodic theorem, Oseledets). *Let $A : \Omega \rightarrow \mathbb{R}^{d \times d}$ be a random matrix on $(\Omega, \mathcal{F}, \mathbb{P}, \varphi)$ and*

$$A_n := \begin{cases} (A \circ \varphi^{n-1}) (A \circ \varphi^{n-2}) \cdots (A \circ \varphi) A, & n \in \mathbb{N}, \\ I, & n = 0, \end{cases}$$

the cocycle on \mathbb{R}^d generated by this sequence. Assume

$$\log^+ \|A\| \in L^1(\Omega, \mathcal{F}, \mathbb{P}).$$

Then there exists $\tilde{\Omega} \in \mathcal{F}$ with $\tilde{\Omega} \subset \varphi^{-1}(\tilde{\Omega})$ and $\mathbb{P}(\tilde{\Omega}) = 1$, such that for each $\omega \in \tilde{\Omega}$ we have:

i) There exists

$$\Psi(\omega) := \lim_{n \rightarrow \infty} (A_n^*(\omega) A_n(\omega))^{1/2n} \geq 0.$$

ii) If

$$e^{\lambda_{p(\omega)}(\omega)} < \dots < e^{\lambda_1(\omega)}$$

are the different eigenvalues of $\Psi(\omega)$ (where $\lambda_{p(\omega)}(\omega) = -\infty$ is possible), and $U_{p(\omega)}(\omega), \dots, U_1(\omega)$ are the corresponding eigenspaces with $d_i(\omega) := \dim U_i(\omega)$, then we have:

$$(\lambda_i \circ \varphi)(\omega) = \lambda_i(\omega), \quad (d_i \circ \varphi)(\omega) = d_i(\omega), \quad \text{and} \quad 1 \leq i \leq p_i(\omega) = (p_i \circ \varphi)(\omega).$$

iii) Defining

$$V_i(\omega) := \begin{cases} \{0\}, & i = p(\omega) + 1 \\ U_{p(\omega)}(\omega) \oplus \dots \oplus U_i(\omega), & i = 1, \dots, p(\omega), \end{cases}$$

we have

$$V_{p(\omega)+1}(\omega) \subset V_{p(\omega)}(\omega) \subset V_{p(\omega)-1}(\omega) \subset \dots \subset V_1(\omega) = \mathbb{R}^d$$

and for each $x \in \mathbb{R}^d \setminus \{0\}$ there exists

$$\lambda(\omega, x) := \lim_{n \rightarrow \infty} \frac{1}{n} \log |A_n(\omega)x|;$$

and we have for all $i = 1, \dots, p(\omega)$:

$$x \in V_i(\omega) \setminus V_{i+1}(\omega) \iff \lambda(\omega, x) = \lambda_i(\omega)$$

resp. equivalently:

$$V_i(\omega) = \{x \in \mathbb{R}^d : \lambda(\omega, x) \leq \lambda_i(\omega)\}.$$

iv) If φ is ergodic, then p, λ_i and d_i on $\tilde{\Omega}$ are constant \mathbb{P} -a.s..

PROOF. By the integrability hypothesis Lemma 7.5 is applicable with $X := \log \|A\|$ and provides the invariant set

$$\tilde{\Omega}_1 := \left\{ \omega \in \Omega : \limsup_{n \rightarrow \infty} \frac{1}{n} \log \|A(\varphi^{n-1}\omega)\| \leq 0 \right\}$$

with full measure. We now apply the deterministic MET Proposition 7.3 to

$$A_n^\omega := A(\varphi^{n-1}\omega) \quad \text{and} \quad \Phi_n^\omega \equiv A_n^\omega \cdots A_1^\omega \stackrel{\text{cocycle}}{=} A_n(\omega),$$

where (4) is valid by definition on $\tilde{\Omega}_1$ and (5) holds true by the theorem of Furstenberg-Kesten 6.15 on a forward invariant set $\tilde{\Omega}_2$ with full measure; consequently Proposition 7.3 is applicable for each $\omega \in \tilde{\Omega}_1 \cap \tilde{\Omega}_2 =: \tilde{\Omega}$, a forward invariant set of full measure, and yields with Theorem 6.15 the desired claims. \square

Definition 7.7. The functions λ_i from the theorem of Oseledets are called *Lyapunov exponents* of the linear cocycle $(A_n)_{n \in \mathbb{N}_0}$.

The spaces V_i (for $i = 1, \dots, p$) are not the analogues of eigenspaces from the deterministic theory. For such an analogy the theory has to be extended to cocycles indexed by \mathbb{Z} , see Arnold [AR 98], Theorem 3.4.11. .

Notations

\mathbb{R}_+	$\{t \in \mathbb{R} : t \geq 0\}$
\mathbb{N}_0	$\mathbb{N} \cup \{0\}$
s^\pm	$(\pm s) \vee 0$; positive resp. negative part of a real number or function s
\equiv	equality by definition
$ \cdot $	norm
$\ \cdot\ $	operator norm
$M \dot{\cup} N$	disjoint union of M and N
$\mathcal{B}(X)$	Borel σ -algebra on the topological space X
\mathcal{B}^n	$\mathcal{B}(\mathbb{R}^n)$
$\delta_1(A) \geq \dots \geq \delta_d(A)$	singular values of $A \in \mathbb{R}^{d \times d}$
$\mathbb{E}(f)$	$\int f d\mathbb{P}$; expectation of a function f with respect to the probability measure \mathbb{P}
$\mathbb{E}(f \mathcal{F})$	conditional expectation of the random variable f given \mathcal{F}
\mathcal{I}	σ -algebra of measurable invariant sets
$\sigma(\mathcal{M})$	σ -Algebra generated by a family \mathcal{M} of sets resp. functions
RV	random variable
Wlog	without loss of generality

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