Decay of correlations for non Hölderian dynamics.  
A coupling approach

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Keywords: dynamical systems, non-Hölder dynamics, mixing rate, chains with complete connections, relaxation speed, coupling methods

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Abstract

We present an upper bound on the mixing rate of the equilibrium state of a dynamical system defined by the one-sided shift and a non Hölder potential of summable variations. The bound follows from an estimation of the relaxation speed of chains with complete connections with summable decay, which is obtained via a explicit coupling between pairs of chains with different histories.

1 Introduction

We consider a dynamical system \((X, T, \mu_\phi)\) where \(X\) is the space of sequences of a finite alphabet, \(T\) is the one-sided shift and \(\mu_\phi\) is the equilibrium state associated to a continuous function \(\phi\) (see Section 2 for precise definitions). We address the question of the speed of convergence of the limit

\[
\int_X f \circ T^n g \, d\mu_\phi \rightarrow_{n \to \infty} \int_X f \, d\mu_\phi \int_X g \, d\mu_\phi.
\]

The dynamical system \((X, T, \mu)\) is said strongly mixing if the convergence occurs for all functions \(f, g\) in a dense subset of \(L^2(\mu_\phi)\). The speed of this convergence —called speed of decay of correlations, or mixing rate— is an important element in the description of the system. In particular, (1.1) determines how averages with respect to measures \(T^n \mu_g\) converge to averages with respect to \(\mu_\phi\), where \(\mu_g = g \mu_\phi / \text{Norm}\).

In this paper we obtain upper bounds showing that for \(\phi\) with summable variations the mixing rate is (at least) summable, polynomial or exponential according to the decay rate of the variations of \(\phi\). The bounds apply for \(f \in L^1(\mu_\phi)\) that do not depend on the future tail-field and \(g\) with variations decreasing proportionally to those of \(\phi\). To obtain these results, we write the difference involved in (1.1) in terms of a chain with complete connections [Doeblin and Fortet (1937), Lalley (1986)] —see identity (3.7) in Section 3 below. The mixing rate is related to the speed with which such chain looses its memory. We bound the latter using coupling techniques.

Previous approaches to the study of the mixing properties of the one-sided shift rely on the use of the transfer operator \(L_\phi\), defined by the duality,

\[
\int_X f \circ T^n g \, d\mu_\phi = \int_X f L^n_\phi g \, d\mu_\phi.
\]

If \(\phi\) is Hölder, this operator, acting on the subspace of Hölder observables, has a spectral gap and the limit (1.1) is attained at exponential speed (Bowen, 1975). When \(\phi\) is not Hölder, the spectral

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gap of the transfer operator may vanish and the spectral study becomes rather complicated. To estimate the mixing rate, Kondah, Maume and Schmitt (1996) proved first that the operator is contracting in the Birkhoff projective metric, while Pollicott (1997), following Liverani (1995), considered the transfer operator composed with conditional expectations. In contrast, our approach is based on a probabilistic interpretation of the duality (1.2) in terms of expectations, conditioned with respect to the past, of a chain with complete connections. The convergence (1.1) is related to the relaxation properties of this chain, which we study via a coupling method.

Chains with complete connections are processes characterized by having transition probabilities that depend on the whole past in a continuous manner. They were first introduced by Onicescu and Mihoc (1935, 1935a) and soon taken up by Doeblin and Fortet (1937). These authors proved the first existence and convergence results for these processes, later extended by Harris (1955). [The definition adopted in these works is written in a manner that differs slightly from current usage in the random-processes literature (see eg. Lalley, 1986). We adopt the latter.] Moreover, their studies were geared towards more complicated objects —called random systems with complete connections— where the chain acts as an underlying “index sequence” used to define very general Markov processes. In this form, the chains have been applied to studies of urn schemes (Onicescu and Mihoc, 1935a), continued fractions (Doeblin 1940; Iosifescu 1978), learning processes (Karlin, 1953; Iosifescu and Theodorescu, 1969; Norman, 1972) and image coding (Barnsley et al, 1988).

As a general reference on the subject we mention the book by Iosifescu (1992) as well as the historical review presented in Kajser (1981) and the brief and clear update by Kajser (1994). These last two references were our main sources for the preceding account. Our work introduces a novel application of this useful objects to the field of dynamical systems, where they appear in a rather natural way.

Coupling ideas were first introduced by Doeblin in his work on the convergence to equilibrium of Markov chains. He let two independent trajectories evolve simultaneously, one starting from the stationary measure and the other from an arbitrary distribution. The convergence follows from the fact that both realizations meet at a finite time. Doeblin wrote his results in 1936, but published them only much later (Doeblin, 1938). [For a description of Doeblin’s contributions to probability theory we refer the reader to Lindvall (1991).] Subsequently, he and Fortet applied this idea to study the existence and relaxation properties of chains with complete connections (Doeblin and Fortet, 1937; see also the account by Iosifescu, 1992). Doeblin’s results can be improved if, instead of letting the trajectories evolve independently, one couples them from the beginning so to reduce the “meeting time” and to ensure that the trajectories evolve together once they meet. Such a procedure is nowadays known as coupling in the stochastic-processes literature. In this setting it is specially efficient to use couplings that “load” the diagonal as much as possible. In our work, we apply a particular coupling with this property, sometimes called the Vaserstein coupling (eg. in Kajser, 1981; Lindval, in his 1992 lectures, calls it γ-coupling). For instance, this coupling prescription applied to a Markov process leads to the so-called Dobrushin’s ergodic coefficient. The sharpness of the convergence rates provided by different types of Markovian couplings has been recently discussed by Burdzy and Kendall (1998).

The Vaserstein coupling has been applied before to chains with complete connections —for instance by Harris, 1955; Kajser (1981, 1994), Barbee (1987) and by us (Bressaud, Fernández, Galves, 1997)— for different purposes. In this paper we use it to obtain a estimation of the relaxation properties of such chains that strengthens and generalizes that of Doeblin and Fortet (1937), even in the form given by Iosifescu (1992). As the transition probabilities of the chains considered here depend on the whole past, the coupling can not ensure that two different trajectories will remain equal after their first meeting time. But the coupling has the property that if the trajectories meet they have a large probability of remaining equal, and this probability increases with the number of consecutive agreements. In the summable case, the coupling is such that with probability one the trajectories disagree only a finite number of times. In fact, the approach can also be applied under an assumption weaker than summability [(4.7) below]. This assumption, which was previously put forward by Harris (1955), ensures that trajectories that differ infinitely often do so with a probability of disagreement that goes to zero. The method leads, in particular,
to a criterion of uniqueness for g-measures proven by Berbee (1987). The mean time between successive disagreements provides a bound on the speed of relaxation of the chain and hence, through the probabilistic interpretation of (1.2), on the mixing rate.

Let us mention, as related developments in the context of dynamical systems, the recent papers by Coelho and Collet (1995) and Young (1997). These papers consider the time two independent systems take to become close. This is reminiscent of the coupling ideas.

The paper is organized as follows. The main results and definitions relevant to dynamical systems are stated in Section 2. The relation between chains with complete connections and the transfer operator is spelled out in Section 3. In Section 4, we state and prove the central result on relaxation speeds of chains with complete connections. Theorem 1 on mixing rates for normalized functions is proven in Section 5, while Theorem 2 on rates for the general case is proven in Section 6. The upper bounds on the decay of correlations depend crucially on estimations of the probability of return to the origin of an auxiliary Markov chain, which are presented in Appendix A.

2 Definitions and statement of the results

Let \( A \) be a finite set henceforth called alphabet. Let us denote
\[
A = \{x = (x_j)_{j \leq -1}, x \in A\}
\]
(2.1)
the set of sequences of elements of the alphabet indexed by the strictly negative integers. Each sequence \( x \in A \) will be called a history. Given two histories \( x \) and \( y \), the notation \( x \equiv y \) indicates that \( x_j = y_j \) for all \( -m \leq j \leq -1 \).

As usual, we endow the set \( A \) with the product topology and the \( \sigma \)-algebra generated by the cylinder sets. We denote by \( C(A, \mathbb{R}) \) the space of real-valued continuous functions on \( A \).

Given a function \( \phi \) on \( A, \phi : A \to \mathbb{R} \), we define its sequence of variations \( (\text{var}_m(\phi))_{m \in \mathbb{N}} \),
\[
\text{var}_m(\phi) = \sup_{x \equiv y} |\phi(x) - \phi(y)| .
\]
(2.2)
We shall say that it has summable variations if,
\[
\sum_{m \geq 1} \text{var}_m(\phi) < +\infty ,
\]
(2.3)
and that it is normalized if it satisfies,
\[
\forall x \in A, \sum_{a \in A} e^{\phi(xa)} = 1 .
\]
(2.4)
We say that a shift-invariant measure \( \mu_\phi \) on \( A \) is compatible with the normalized function \( \phi \) if and only if, for \( \mu_\phi \)-almost-all \( x \) in \( A \),
\[
E_{\mu_\phi} \left( 1_{\{x_{-1} = a\}} | \mathcal{F}_{\leq -2} \right) (x) = e^{\phi(T(x)a)} ,
\]
(2.5)
where the left-hand side is the usual conditional expectation of the the indicator function of the event \( \{x_{-1} = a\} \) with respect to the \( \sigma \)-algebra of the past up to time \(-2\).

An equivalent way of expressing this is by saying that \( \mu_\phi \) is a g-measure for \( g = e^\phi \). If \( \phi \) has summable variations, and even under a slightly weaker conditions, then such a measure is unique and will be denoted \( \mu_\phi \). The measure \( \mu_\phi \) can also be characterized via a variational principle, in which context it is called equilibrium state for \( \phi \). For details see Ledrappier (1974), Walters (1975), Quas (1996) and Berbee (1987).
For a non-constant \( \phi \), we consider the seminorm
\[
\|g\|_\phi = \sup_{k \geq 0} \frac{\text{var}_k(g)}{\text{var}_k(\phi)}
\] (2.6)
and the subspace of \( C(A, \mathbb{R}) \) defined by,
\[
V_\phi = \left\{ g \in C(A, \mathbb{R}) , \|g\|_\phi < +\infty \right\}. \tag{2.7}
\]

Given a real-valued sequence \((\gamma_n)_{n \in \mathbb{N}}\), let \((S_n^{(\gamma)})_{n \in \mathbb{N}}\) be the Markov chain taking values in the set \( \mathbb{N} \) of natural numbers starting from the origin
\[
\mathbf{P}(S_0^{(\gamma)} = 0) = 1 \tag{2.8}
\]
whose transition probabilities are defined by
\[
p_{i,i+1} = 1 - \gamma_i, \quad p_{i,0} = \gamma_i, \tag{2.9}
\]
for all \( i \in \mathbb{N} \). For any \( n \geq 1 \) we define
\[
\gamma_n^* = \mathbf{P}(S_n^{(\gamma)} = 0). \tag{2.10}
\]

We now state our first result.

**Theorem 1** Let \( \phi : A \to \mathbb{R} \) be a normalized function with summable variations and set
\[
\gamma_n = 1 - e^{-\text{var}_n(\phi)}. \tag{2.11}
\]
Then,
\[
\left| \int f \circ T^n g \, d\mu_\phi - \int g \, d\mu_\phi \right| \leq \|f\|_1 \|g\|_\phi \sum_{k=0}^{n} \text{var}_k(\phi) \gamma_n^* - k \tag{2.12}
\]
\[
\leq C \|f\|_1 \|g\|_\phi \gamma_n^*, \tag{2.13}
\]
for all \( g \in V_\phi \) and \( f \in L^1(\mu_\phi) \) measurable with respect to \( \bigcup_{i \in \mathbb{Z}} \mathcal{F}_{\leq i} \). The constant \( C \) can be explicitly computed.

This theorem is proven in Section 5, using the results obtained in Section 4 on the relaxation speed of chains with complete connections.

For each non-normalized function \( \phi \) with summable variations there exist a unique positive function \( \rho \) and a unique real number \( c \) such that the function
\[
\psi = \phi + \log \rho - \log \rho \circ T + c \tag{2.14}
\]
is normalized (Walters, 1975). We call \( \psi \) the normalization of \( \phi \). The construction of compatible measures given in (2.5) loses its meaning for non-normalized \( \phi \). It is necessary to resort to an alternative characterization in terms of a variational principle (see eg. Bowen 1975) leading to equilibrium states. In Walters (1975) it is proven that:

(a) \( \phi \) with summable variations admits a unique equilibrium state, that we denote also \( \mu_\phi \);

(b) the corresponding normalized \( \psi \), given by (2.14), admits a unique compatible measure \( \mu_\psi \)
    (even when the variations of \( \psi \) may not be summable), and

(c) \( \mu_\phi = \mu_\psi \).
Our second theorem generalizes Theorem 1 to non-normalized functions.

**Theorem 2** Let $\phi : A \to \mathbb{R}$ be a function with summable variations and let $\psi$ be its normalization. Let $(n_m)_{m \in \mathbb{N}}$ be an increasing subadditive sequence such that the subsequence of the rests, $(\sum_{k \geq n_m} \text{var}_k(\phi))_{m \geq 0}$, is summable, and

$$\gamma_m = 1 - e^{-3 \sum_{k \geq n_m} \text{var}_k(\phi)} ;$$  \hspace{1cm} (2.15)

then,

$$\left| \int f \circ T^n g \, d\mu_{\phi} - \int f \, d\mu_{\phi} \int g \, d\mu_{\phi} \right| \leq ||f||_1 ||g||_1 \sum_{k=0}^{n} \text{var}_{k}(\phi) \gamma_{n-k} $$  \hspace{1cm} (2.16)

$$\leq C ||f||_1 ||g||_1 \gamma_n ,$$  \hspace{1cm} (2.17)

for all $g \in V_{\phi}$ and $f \in L^1(\mu_{\phi})$ measurable with respect to $\bigcup_{i \in \mathbb{Z}} F_{\leq i}$, where $C$ is a computable constant. Here $\gamma_n$ is defined as in (2.10) but using the sequence $(\gamma_n)_{n \in \mathbb{N}}$.

The estimation of the large-$n$ behavior of the sequence $(\gamma_n)_{n \in \mathbb{N}}$ given the behavior of the original $(\gamma_n)_{n \in \mathbb{N}}$ only requires elementary computations. For the convenience of the reader we summarize some results in Appendix A.

### 3 Transfer operators and chains.

Let $P$ be a family of transition probabilities on $A \times A$.

$$P : \begin{cases} A \times A & \mapsto [0; 1] \\ (a, z) & \mapsto P(a \mid z) . \end{cases}$$  \hspace{1cm} (3.1)

Given a history $x$, a chain with past $x$ and transitions $P$, is the process $(Z^n_x)_{n \in \mathbb{N}}$ whose conditional probabilities satisfy

$$P(Z^n_x = a \mid Z^n_{x+j} = z_j, j \leq -1) = P(a \mid z) \text{ for } n \geq 0 ,$$  \hspace{1cm} (3.2)

for all $a \in A$ and all histories $z$ with $z_{j-n} = x_j, j \leq -1$, and such that

$$Z^n_x = x_n , \text{ for } n \leq -1 .$$  \hspace{1cm} (3.3)

This chain can be interpreted as a conditioned version of the process defined by the transition probabilities (3.1), given a past $x$ (for more details, see Quas 1996).

Let $\phi : A \to \mathbb{R}$ be a continuous normalized function. The transfer operator associated to $\phi$ is the operator $L_\phi$ acting on $C(A, \mathbb{R})$ defined by,

$$L_\phi f(x) = \sum_{y : T(y) = x} e^{\phi(y)} f(y) .$$  \hspace{1cm} (3.4)

This operator is related to the conditional probability (2.5) in the form

$$E_{\mu_{\phi}}(f \mid F_{\leq -2}) = (L_\phi f) \circ T .$$  \hspace{1cm} (3.5)

This relation shows the equivalence of (1.2) and (3.4) as definitions of the operator. In addition, if $\phi$ is normalized we can construct, for each history $x \in A$, the chain $Z^n_x(\phi) = (Z^n_x)_{n \in \mathbb{Z}}$ with past $x$ and transition probabilities

$$P(a \mid x) = e^{\phi(xa)} .$$  \hspace{1cm} (3.6)
Iterates of the transfer operator, $L^n \phi g(x)$, on functions $g \in C(A, \mathbb{R})$ can be interpreted as expectations $E[g((Z_{n+j}^x)_{j \leq -1})]$ of the chain. Indeed,

$$L^n \phi g(x) = \sum_{a_1, \ldots, a_n \in A} e^{\sum_{k=1}^n \phi(xa_1 \cdots a_k)} g(xa_1 \cdots a_n)$$

$$= \sum_{a_1, \ldots, a_n \in A} \left( \prod_{k=1}^n P(a_k | a_{k-1} \cdots a_1 x) \right) g(xa_1 \cdots a_n)$$

$$= E[g((Z_{n+j}^x)_{j \leq -1})].$$

From this expression and the classical duality (1.2) between the composition by the shift and the transfer operator $L \phi$ in $L^2(\mu \phi)$, we obtain the following expression for the decay of correlations,

$$\int f \circ T^n g \, d\mu \phi - \int f \, d\mu \phi \int g \, d\mu \phi$$

$$= \int f(x) L^n \phi g(x) \, d\mu \phi(x) - \int f(x) \left( \int L^n \phi g(y) \, d\mu \phi(y) \right) \, d\mu \phi(x)$$

$$= \int f(x) \left( E[g((Z_{n+j}^x)_{j \leq -1})] - E[g((Z_{n+j}^y)_{j \leq -1})] \right) \, d\mu \phi(y) \, d\mu \phi(x). \quad (3.7)$$

This inequality shows how the speed of decay of correlations can be bounded by the speed with which the chain looses its memory. We deal with the later problem in the next section.

4 Relaxation speed for chains with complete connections

4.1 Definitions and main result

We consider chains whose transition probabilities satisfy

$$\inf_{x, y: x \equiv y} \frac{P(a | x)}{P(a | y)} \geq 1 - \gamma_m,$$ \quad (4.1)

for some real-valued sequence $(\gamma_m)_{m \in \mathbb{N}}$, decreasing to 0 as $m$ tends to $+\infty$. Without loss of generality, this decrease can be assumed to be monotonic. To avoid trivialities we assume $\gamma_0 < 1$. In the literature, a stationary process satisfying (4.1) is called a chain with complete connections.

For a set of transition probabilities satisfying (4.1), we consider, for each $x \in A$, the chain $(Z_{n+j}^x)_{n \in \mathbb{Z}}$ with past $x$ and transitions $P$ [see (3.2)–(3.3)]. The following proposition plays a central role in the proof of our results.

**Proposition 1** For all histories $x, y \in A$, there is a coupling $((\tilde{U}_{n,y}^x, \tilde{V}_{n,y}^x))_{n \in \mathbb{Z}}$ of $(Z_{n+j}^x)_{n \in \mathbb{Z}}$ and $(Z_{n+j}^y)_{n \in \mathbb{Z}}$ such that the integer-valued process $(T_{n,y}^x)_{n \in \mathbb{Z}}$ defined by

$$T_{n,y}^x = \inf\{m \geq 0 : \tilde{U}_{n-m}^y = \tilde{V}_{n-m}^y \}, \quad (4.2)$$

satisfies

$$P(T_{n,y}^x = 0) \leq \gamma_n^*, \quad (4.3)$$

for $n \geq 0$, where $\gamma_n^*$ was defined in (2.10).

The proof of this proposition is given in Section 4.4.
An immediate consequence of this proposition is the following bound on the relaxation rate of the processes $Z^x$.

**Corollary 1** For all histories $x$ and $y$, for all $a \in A$,

$$|P(Z^x_n = a) - P(Z^y_n = a)| \leq \gamma^*_n,$$

(4.4)

and, for $k \geq 1$,

$$|P(Z^x_n, \ldots, Z^x_{n+k}) = (a_0, \ldots, a_k) - P(Z^y_n, \ldots, Z^y_{n+k}) = (a_0, \ldots, a_k)|$$

$$\leq \sum_{j=0}^{k} \left( \prod_{m=1}^{j-1} (1 - \gamma_m) \right) \gamma^*_{n+k-j}.$$  

(4.5)

This corollary is proved in Section 4.5.

**Remark 1** Whenever $\gamma^*_n \to 0$,

(4.6)

inequality (4.4) implies the existence and uniqueness of the invariant measure compatible with a system of conditional probabilities satisfying (4.1). In fact, property (4.6) holds under the condition

$$\sum_{m \geq 1} k \prod_{m=1}^{k} (1 - \gamma_k) = +\infty,$$

(4.7)

which is weaker than summability. In this case, the Markov chain $(S_n)_{n \in \mathbb{N}}$ is no longer transient but it is null recurrent and the property $P(S_n = 0) \to 0$ remains true.

**Remark 2** If $X = (X_n)_{n \in \mathbb{Z}}$ is a stationary process with transition $P$ satisfying (4.1), then Corollary 1 implies

$$|P(Z^x_n = a) - P(X_n = a)| \leq \gamma^*_n,$$

(4.8)

uniformly in the history $x$.

**Remark 3** Inequality (4.3) constitutes a double improvement of the Doeblin-Fortet-Iosifescu results (see Theorem 1 in Iosifescu, 1992). First, the latter only apply when the remainders $\sum_{i \geq n} \gamma_i$ tend to zero with $n$. Second, even for this case, Doeblin-Fortet-Iosifescu estimation can be seen to be an upper bound for our $\gamma^*_n$. The bound is in general strictly larger.

### 4.2 Maximal coupling

Given two probability distributions $\mu = (\mu(a))_{a \in A}$ and $\nu = (\nu(a))_{a \in A}$ we denote by $\mu \tilde{\times} \nu = (\mu \tilde{\times} \nu(a, b))_{(a, b) \in A \times A}$ the so-called maximal coupling of the distributions $\mu$ and $\nu$ defined as follows:

$$\begin{cases} 
\mu \tilde{\times} \nu(a, a) = \mu(a) \wedge \nu(a) & \text{if } a = b \\
\mu \tilde{\times} \nu(a, b) = \frac{(\mu(a) - \nu(a))^{+}(\nu(b) - \mu(b))^{+}}{\sum_{e \in A} (\mu(e) - \nu(e))^{+}} & \text{if } a \neq b.
\end{cases}$$

(4.9)

For more details on maximal couplings see Appendix A.1 in Barbour, Holst and Janson (1992).

The coupling is maximal in the sense that the distribution $\mu \tilde{\times} \nu$ on $A \times A$ maximizes the weight

$$\Delta(\zeta) = \sum_{a \in A} \zeta(a,a)$$

of the diagonal among the distributions $\zeta$ on $A \times A$ satisfying simultaneously

$$\sum_{a \in A} \zeta(a,b) = \nu(b) \quad \text{and} \quad \sum_{b \in A} \zeta(a,b) = \mu(a).$$
For this coupling, the weight $\Delta(\mu \times \nu)$ of the diagonal satisfies,

$$\Delta(\mu \times \nu) = \sum_{a \in A} \mu(a) \land \nu(a) = 1 - \sum_{a \in A} (\mu(a) - \nu(a))^+ = 1 - \frac{1}{2} \sum_{a \in A} |\mu(a) - \nu(a)|. \quad (4.10)$$

Moreover,

$$\Delta(\mu \times \nu) = 1 - \sum_{a \in A} \mu(a) \left( 1 - \frac{\nu(a)}{\mu(a)} \right)^+ \geq 1 - \sum_{a \in A} \mu(a) \left( 1 - \inf_{a' \in A} \frac{\nu(a')}{\mu(a')} \right) = \inf_{a \in A} \frac{\nu(a)}{\mu(a)} \quad (4.11)$$

### 4.3 Coupling of chains with different pasts

Given a double history $(x, y)$, we consider the transition probabilities defined by the maximal coupling

$$\tilde{P}((a, b) | x, y) = [P(\cdot | x) \tilde{x} P(\cdot | y)] (a, b). \quad (4.12)$$

By (4.1) we have,

$$\inf_{a \in A, u \equiv v} \frac{P(a | u)}{P(a | v)} \geq 1 - \gamma_m.$$ 

By (4.11) this implies that

$$\Delta\left(\tilde{P}(\cdot, \cdot | x, y)\right) \geq 1 - \gamma_m, \quad (4.13)$$

everywhere $x \equiv y$.

Now, we fix a double history $(x, y)$ and we define $(\tilde{U}_n^{x,y}, \tilde{V}_n^{x,y})_{n \in \mathbb{Z}}$ to be the chain taking values in $A^2$, with past $(x, y)$ and transition probabilities given by (4.12). If $x \equiv y$, (4.13) yields

$$\mathbb{P}(\tilde{U}_0^{x,y} \neq \tilde{V}_0^{x,y}) \leq \gamma_m. \quad (4.14)$$

We denote

$$\Delta_{m, n} := \left\{ \tilde{U}_j = \tilde{V}_j, m \leq j \leq n \right\}. \quad (4.15)$$

Notice that $\Delta_{-m, -1}$ is the union over all the sequences $x, y$ with $x \equiv y$ of the events $\{(\tilde{U}_j, \tilde{V}_j) = (x_j, y_j) ; j \leq -1\}$. Using the stationarity of the conditional probabilities, we obtain

$$\mathbb{P}(\tilde{U}_n \neq \tilde{V}_n | \Delta_{n-m, n-1}) \leq \gamma_m, \quad (4.16)$$

for all $n \geq 0$.

### 4.4 Proof of Proposition 1

From this subsection on, we will be working with bounds which are uniform in $x, y$, hence we will omit, with a few exceptions, the superscript $x, y$ in the processes $T_n^{x,y}$ (defined below), $\tilde{U}_n^{x,y}$ and $\tilde{V}_n^{x,y}$.

Let us consider the integer-valued process $(T_n)_{n \in \mathbb{Z}}$ defined by:

$$T_n = \inf\{m \geq 0 : \tilde{U}_{n-m} \neq \tilde{V}_{n-m}\}. \quad (4.17)$$

For each time $n$, the random variable $T_n$ counts the number of steps backwards needed to find a difference in the coupling. First, notice that (4.16) implies that,

$$\mathbb{P}(T_{n+1} = k + 1 | T_n = k) \geq 1 - \gamma_k \quad (4.18)$$

and

$$\mathbb{P}(T_{n+1} = 0 | T_n = k) \leq \gamma_k, \quad (4.19)$$

all the other transition probabilities being zero. This process $(T_n)_{n \in \mathbb{Z}}$ is not a Markov chain.
We now consider the integer-valued Markov chain \((S_n^{(\gamma)})_{n \geq 0}\) starting from state 0 and with transition probabilities given by (2.9), that is \(p_{i,i+1} = 1 - \gamma_i \) and \(p_{i,0} = \gamma_i\). Proposition 1 follows from the following lemma, setting \(k = 1\).

**Lemma 1** For each \(k \in \mathbb{N}\), the following inequality holds

\[
P(S_n^{(\gamma)} \geq k) \leq P(T_n \geq k)
\]  

(4.20)

**Proof** We shall proceed by induction on \(n\). Since \(P(S_0^{(\gamma)} = 0) = 1\), inequalities (4.20) hold for \(n = 0\). Assume now that (4.20) holds for some integer \(n\). There is nothing to prove for \(k = 0\). For \(k \geq 1\),

\[
P(T_{n+1} \geq k) = \sum_{m=k}^{+\infty} P(T_{n+1} = m)
= \sum_{m=k}^{+\infty} P(T_{n+1} = m \mid T_n = m - 1) P(T_n = m - 1)
\geq \sum_{m=k}^{+\infty} (1 - \gamma_{m-1}) P(T_n = m - 1)
= \sum_{m=k}^{+\infty} (1 - \gamma_{m-1}) \left( P(T_n \geq m - 1) - P(T_n \geq m) \right)
= (1 - \gamma_k) P(T_n \geq k - 1) + \sum_{m=k}^{+\infty} (\gamma_{m-1} - \gamma_m) P(T_n \geq m).
\]  

(4.21)

By the same computation, we see that

\[
P(S_{n+1}^{(\gamma)} \geq k) = (1 - \gamma_k) P(S_n^{(\gamma)} \geq k - 1) + \sum_{m=k}^{+\infty} (\gamma_{m-1} - \gamma_m) P(S_n^{(\gamma)} \geq m).
\]  

(4.22)

Hence, using the recurrence assumption and the fact that \((\gamma_n)_{n \geq 0}\) is decreasing we conclude that

\[
P(T_{n+1} \geq k) \geq P(S_{n+1}^{(\gamma)} \geq k),
\]

for all \(k \geq 1\). \(\Box\)

### 4.5 Proof of Corollary 1

To prove (4.4), first notice that by construction the process \((\tilde{U}_n)_{n \in \mathbb{Z}}\) has the same law as \((Z_n^x)_{n \in \mathbb{Z}}\) and \((\tilde{V}_n)_{n \in \mathbb{Z}}\) has the same law as \((Z_n^y)_{n \in \mathbb{Z}}\). Thus,

\[
\left| P(Z_n^x = a) - P(Z_n^y = a) \right| = \left| P(\tilde{U}_n = a) - P(\tilde{V}_n = a) \right| \leq P(\tilde{U}_n \neq \tilde{V}_n)
\]  

(4.23)

Hence, by definition of the process \(T_n\) and Lemma 1,

\[
\left| P(Z_n^x = a) - P(Z_n^y = a) \right| \leq P(T_n = 0) \leq P(S_n^{(\gamma)} = 0).
\]  

(4.24)

The proof of (4.5) starts similarly:

\[
\left| P((Z_n^x, \ldots, Z_{n+k}^x) = (a_0, \ldots, a_k)) - P((Z_n^y, \ldots, Z_{n+k}^y) = (a_0, \ldots, a_k)) \right|
\leq P(S_{n+k}^{(\gamma)} \leq k + 1).
\]

To conclude, we notice that

\[
P(S_n^{(\gamma)} \leq k) = \sum_{j=0}^{k} P(S_n^{(\gamma)} = j) = \sum_{j=0}^{k} \left( \prod_{m=1}^{j-1} (1 - \gamma_m) \right) P(S_{n-j}^{(\gamma)} = 0).
\]  

\(\Box\)  

(4.25)
5 Proof of Theorem 1

The proof of Theorem 1 is based on the inequality

$$\left| \int f \circ T^x \, g \, d\mu - \int f \, d\mu \int g \, d\mu \right| \leq ||f||_1 \sup_{x,y} E \left[ g((\tilde{U}^{x,y}_{n+j})_{j \leq -1}) - g((\tilde{V}^{x,y}_{n+j})_{j \leq -1}) \right],$$  

(5.1)

which follows from (3.7) and the fact that $((\tilde{U}^{x,y}_n, \tilde{V}^{x,y}_n))_{n \in \mathbb{Z}}$ is a coupling between the chains with pasts $x$ and $y$, respectively. An upper bound to the right-hand side is provided by Proposition 1. We see that the transition probabilities (3.6) satisfy condition (4.1), since

$$\frac{P(a \mid x)}{P(a \mid y)} = e^{\varphi(ax) - \varphi(ay)} \geq e^{-\text{var}_{m+1}(\varphi)}$$  

(5.2)

whenever $x, y \in \mathbb{A}$ are such that $x = \gamma y$ for some $m \in \mathbb{N}$. We can therefore apply Proposition 1 with

$$\gamma_m = 1 - e^{-\text{var}_{m+1}(\varphi)},$$  

(5.3)

which tends monotonically to zero if $\sum_{m \geq 1} \text{var}_m(\varphi) < +\infty$.

To prove (2.12) we use the process $(T^{x,y}_n)_{n \in \mathbb{N}}$ to obtain the upper bound

$$E \left[ g((\tilde{U}^{x,y}_{n+j})_{j \leq -1}) - g((\tilde{V}^{x,y}_{n+j})_{j \leq -1}) \right] = E \left[ \sum_{k=0}^{+\infty} \text{var}_k(g) \, P(T^{x,y}_n = k) \right] \leq \sum_{k=0}^{+\infty} \text{var}_k(g) \, P(T^{x,y}_n = k),$$  

(5.4)

Now, in order to use the bound (4.3) of Proposition (1) we resort to the monotonicity of the variations of $\varphi$:

$$\sum_{k=0}^{+\infty} \text{var}_k(\varphi) \, P(T^{x,y}_n = k) \leq \sum_{k=0}^{n-1} \text{var}_k(\varphi) \, P(T^{x,y}_n = k) + \text{var}_n(\varphi) \, \sum_{k=n}^{+\infty} P(T^{x,y}_n = k)$$

$$= \sum_{k=0}^{n-1} \text{var}_k(\varphi) \, P(T^{x,y}_{n-k} = 0) + \text{var}_n(\varphi) \, \sum_{k=n}^{+\infty} P(T^{x,y}_0 = k - n)$$

$$\leq \sum_{k=0}^{n} \text{var}_k(\varphi) \, P(S^{(\gamma)}_{n-k} = 0),$$  

(5.5)

uniformly in $x, y$. The bound (2.12) follows from (5.1), (5.4), (5.5) and the fact that

$$\sum_{j=0}^{+\infty} P(T^{x,y}_0 = j) = 1 = P(S^{(\gamma)}_0 = 0).$$  

(5.6)

To prove (2.13) we use the strong Markov property of the process $(S^{(\gamma)}_n)_{n \in \mathbb{N}}$ to obtain

$$P(S^{(\gamma)}_n = 0) = \sum_{k=1}^{n} P(\tau = k) \, P(S^{(\gamma)}_{n-k} = 0),$$  

(5.7)

where

$$\tau = \inf\{n > 0; S^{(\gamma)}_n = 0\}.$$  

(5.8)
We now use (5.7) to bound the last line in (5.5) in the form

\[
\sum_{k=0}^{\gamma_n} \text{var}_k(\phi) P(S_{n-k}^{(\gamma)} = 0) \leq \sum_{k=1}^{\gamma_n} [\text{var}_0(\phi) P(\tau = k) + \text{var}_k(\phi)] P(S_{n-k}^{(\gamma)} = 0)
\]

\[
\leq C \sum_{k=1}^{\gamma_n} P(\tau = k) P(S_{n-k}^{(\gamma)} = 0)
\]

\[
= C P(S_{\gamma_n}^{(\gamma)} = 0),
\]

with

\[
C = \text{var}_0(\phi) + \sup_k \frac{\text{var}_k(\phi)}{P(\tau = k)}. \tag{5.10}
\]

To conclude, we must prove that the constant C is finite. By direct computation,

\[
P(\tau = 1) = \gamma_0,
\]

\[
P(\tau = n) = \gamma_{n-1} \prod_{m=0}^{n-2} (1 - \gamma_m) \text{ for } n \geq 2, \tag{5.11}
\]

\[
P(\tau = +\infty) = \prod_{m=0}^{+\infty} (1 - \gamma_m).
\]

From this and (2.11) we obtain

\[
\lim_{k \to \infty} \frac{\text{var}_k(\phi)}{P(\tau = k)} = \lim_{k \to \infty} \frac{\text{var}_k(\phi)}{1 - e^{-\text{var}_k(\phi)}} \prod_{m=0}^{k-2} (1 - \gamma_m). \tag{5.12}
\]

Since var\(_k(\phi) \to 0\), the first fraction converges to 1. We see from (5.11) that the second fraction converges to 1/P(\tau = +\infty). By elementary calculus, this is finite since \(\phi\) has summable variations.

**Remark 4** The previous computations lead to stronger results for more regular functions \(g\). For example, when \(g\) satisfies

\[
\text{var}_k(g) \leq ||g||_{\theta} \theta^k
\]

for some \(\theta < 1\) and some \(||g||_{\theta} < \infty\) (Hölder norm of \(g\)), a chain of inequalities almost identical to those ending in (5.4) leads to

\[
\left| \int f \circ T^n g d\mu - \int f d\mu \int g d\mu \right| \leq ||f||_1 \sum_{k=0}^{+\infty} ||g||_{\theta} \theta^k \gamma_{n-k}^* \leq ||f||_1 ||g||_{\theta} \sum_{k=0}^{n} \theta^{-k} \gamma_k^*. \tag{5.14}
\]

On the other hand, if \(g\) is a function that depends only on the first coordinate, we get,

\[
\left| \int f \circ T^n g d\mu - \int f d\mu \int g d\mu \right| \leq ||f||_1 \sup_{x,y} \left| E[g(Z_{n}^x)] - E[g(Z_{n}^y)] \right| \leq ||f||_1 ||g||_{\infty} P(\tilde{U}_n \neq \tilde{V}_n) \leq ||f||_1 ||g||_{\infty} \gamma_n^* . \tag{5.15}
\]

### 6 Proof of Theorem 2

We now consider the general case where the function \(\phi\) is not necessarily normalized. In this case we resort to the normalization \(\psi\) defined in (2.14) and we consider chains with transition probabilities

\[
P(a \mid x) = e^{\phi(xa)} \frac{\rho(xa)}{\rho(x)} e^c =: e^{\psi(xa)}. \tag{6.1}
\]
However, the summability of the variations of $\phi$ does not imply the analogous condition for $\psi$, because there are additional “oscillations” due to the cocycle $\log \rho - \log \rho \circ T$. Instead,

$$\text{var}_m \psi \leq \sum_{k \geq m} \text{var}_k(\phi),$$

for all $m \geq 0$ (see Walters, 1978). Hence, we can apply Theorem 1 only under the condition

$$\sum_{k=1}^{+\infty} k \text{var}_k(\phi) < +\infty.$$  

If this is the case, the correlations for functions $f \in L^1(\mu)$ and $g \in V_\psi$ decay faster than $\gamma^*_m$, where $\gamma_m = e^{\sum_{k \geq m} \text{var}_k(\phi)} - 1$.

To prove the general result without assuming (6.3) we must work with block transition probabilities, which are less sensitive to the oscillations of the cocycle. More precisely, given a family of transition probabilities $P$ on $A \times A$, let $P_n$ denote the corresponding transition probabilities on $A^n \times A^n$:

$$P_{n+1}(a_0,n \mid x) = P(a_n \mid a_{n-1} \cdots a_1 x) \cdots P(a_2 \mid a_1 x) P(a_1 \mid x)$$

where $a_0,n := (a_0, \ldots, a_n) \in A^{n+1}$. If the transition probabilities $P$ are defined by a normalized function $\phi$ as in (3.6), then we see from (6.4) that the transition probabilities $P_n$ obey a similar relation

$$P_n(a_0,n-1 \mid x) = e^{\phi_n(xa_0,n-1)},$$

with

$$\phi_n(xa_0,n-1) := \sum_{k=0}^{n-1} \phi(xa_0 \cdots a_k).$$

In particular, for transitions (6.1) the formula (6.4) yields

$$\psi_n = \phi_n + \log \rho - \log \rho \circ T + nc.$$  

A comparison of (6.8) with (6.2) shows that it is largely advantageous to bound directly the oscillations of $\psi_n$. This is what we do in this section by adapting the arguments of Section 5.

### 6.1 Coupling of the transition probabilities for blocks

For every integer $n$, we define a family of transition probability $\mathcal{P}_n$ on $(A^n)^2 \times A^2$ by

$$\mathcal{P}_n(a_0,n-1, b_0,n-1 \mid x, y) = [P_n(\cdot \mid x) \times P_n(\cdot \mid y)](a_0,n-1; b_0,n-1).$$

Let $(n_m)_{m \in \mathbb{N}}$ be an increasing sequence. For each double history $x, y$, we consider the coupling of the chains for $n_m$-blocks with past $x$ and $y$, defined by,

$$\mathcal{P}(\mathcal{U}^{x,y}_{0,n_m} = a_{0,n_m} , \mathcal{V}^{x,y}_{0,n_m} = b_{0,n_m})$$

$$= \prod_{m=1}^{M} \mathcal{P}_{n_{m+1} - n_m}(a_{n_m,n_{m+1}} , b_{n_m,n_{m+1}} \mid a_{n_m} \cdots a_0x , b_{n_m} \cdots b_0y).$$
6.2 The process of last block-differences

We set
\[ \gamma_{k}^{(n)} = 1 - \inf \left\{ \frac{P_{n}(a_{0},n-1|y)}{P_{n}(a_{0},n-1|x)} : x = y, a_{1}, \ldots, a_{n-1} \in A \right\}. \] (6.11)

From (4.11) we see that, for \( x = y \), the weight of the diagonal of each coupling \( P_{n} \) satisfies
\[ \Delta(P_{n}(\cdot, \cdot|x,y)) \geq \inf_{a_{0},\ldots,a_{n-1} \in A} P_{n}(a_{0},n-1|x)P_{n}(a_{0},n-1|y) \geq 1 - \gamma_{k}^{(n)}. \] (6.12)

If we denote
\[ \Delta_{x,y}^{m,k} := \left\{ U_{i}^{x,y} = V_{i}^{x,y} \right\}_{n_{m} \leq j \leq n_{m+k}}, \]
we deduce from (6.12) that
\[ P(\Delta_{m+k+1}^{x,y} | \Delta_{m,m+k}) \geq 1 - \gamma_{n_{m+k}-n_{m}}^{(n_{m+k+1}-n_{m+k})}. \] (6.13)

We construct the process \((T_{n})_{n \in \mathbb{N}}\) with
\[ T_{m}^{x,y} = \inf \left\{ p \geq 0 : U_{i}^{x,y} \neq V_{i}^{x,y} \text{ for some } i, n_{m-p} \leq i \leq n_{m-p+1} \right\}. \] (6.14)

By (6.13), the conditional laws of this process satisfy,
\[ P(T_{m+1} = k+1 | T_{m} = k) \geq 1 - \gamma_{n_{m+k+1}-n_{m}}^{(n_{m+k+1}-n_{m})}, \] (6.15)
and
\[ P(T_{m+1} = 0 | T_{m} = k) \leq \gamma_{n_{m+k+1}-n_{m}}^{(n_{m+k+1}-n_{m})}. \] (6.16)

6.3 The dominating Markov process

Let us choose the length of the blocks in such a way that the sequence \((n_{m})_{m \in \mathbb{N}}\) is subadditive, i.e.
\[ n_{m+k} - n_{m} \leq n_{k} \] (6.17)
for \( m, k \geq 0 \), and that
\[ \sup_{n \geq 0} \gamma_{n}^{(n)} < 1 \] (6.18)
for all \( \ell \geq 0 \). These two properties together with (6.15)–(6.16) imply that, for all histories \( x \) and \( y \),
\[ P(T_{m+1}^{x,y} = k+1 | T_{m}^{x,y} = k) \geq 1 - \tau_{k}^{(n_{m+k+1}-n_{m})}, \] (6.19)
and
\[ P(T_{m+1}^{x,y} = 0 | T_{m}^{x,y} = k) \leq \tau_{k}^{(n_{m+k+1}-n_{m})}. \] (6.20)
with
\[ \tau_{k} := \sup_{n \geq 1} \gamma_{n}^{(n)}, \] (6.21)
for \( m \geq 1 \).

We now define the “dominating” Markov chain \((S_{n}^{\gamma})_{n \in \mathbb{N}}\) as in (2.8)–(2.9). Lemma 1 yields
\[ P(T_{m}^{x,y} = 0) \leq P(S_{m} = 0) \leq \tau_{m}^{*}. \] (6.22)
Hence, if \( n_{m} \leq n \leq n_{m+1} \),
\[ P(U_{n}^{x,y} \neq V_{n}^{x,y}) \leq P(T_{m}^{x,y} = 0) \leq \tau_{m}^{*}. \] (6.23)

6.4 Decay of correlations

We can now mimic the proof of Theorem 5 in terms of barred objects.
As \( (\text{var}_m(\phi))_{m \in \mathbb{N}} \) is summable, there exists a subadditive sequence \( (n_m)_{m \in \mathbb{N}} \) such that the sequence \( \alpha_m \) of the tails

\[
\alpha_m = \sum_{k \geq n_m} \text{var}_k(\phi)
\]

is summable:

\[
\sum_{m \geq 0} \alpha_m < +\infty.
\]

(6.25)

The transitions for blocks of size \( n \) satisfy

\[
\frac{P_n(a_0,n-1|x)}{P_n(a_0,n-1|y)} \geq e^{-\text{var}_k(\psi_n)}
\]

(6.26)

if \( x = y \). But from (6.8), (6.7) and (6.2) we have

\[
\text{var}_k(\psi_n) \leq \left( \sum_{m=k}^{k+n} \sum_{m=k+n}^{m+k} \sum_{m \geq k} \text{var}_m(\phi) \right)
\]

\[
\leq 3 \sum_{m \geq k} \text{var}_m(\phi).
\]

(6.27)

Hence we can choose in (6.21)

\[
\gamma_k \leq 1 - e^{-3\alpha_k},
\]

(6.28)

a choice for which

\[
\sum_{k \geq 1} \gamma_k < +\infty.
\]

(6.29)

To prove the theorem, we now proceed as in (5.1) and (5.4)–(5.10) but replacing tildes by bars and putting bars over the processes \((T_n)\) and \((S_n^{(\gamma)})\). We just point out that, due to the subadditivity of \( n_m \),

\[
\text{var}_{(n_m+k-n_m)}(\phi) \leq \text{var}_{n_k}(\phi)
\]

uniformly in \( m \).

A Returns to the origin of the dominating Markov chain

In this appendix we collect a few results concerning the probability of return to the origin of the Markov chain \((S_n^{(\gamma)})_{n \in \mathbb{N}}\) defined via (2.9). (In the sequel we omit the superscript “\((\gamma)\)” for simplicity.) We point out that the right-hand-side of the displayed formula in Theorem 1 of Iosifescu (1992) is indeed an upper bound for \( P(S_n = 0) \). This bound is clearly less sharp than our estimations, for instance when \((\gamma_m)\) decreases polynomially.

**Proposition 2** Let \((\gamma_m)_{n \in \mathbb{N}}\) be a real-valued sequence decreasing to 0 as \( n \to +\infty \).

(i) If \( \sum_{m \geq 1} \prod_{k=0}^{m} (1 - \gamma_k) = +\infty \), then \( P(S_n = 0) \to 0 \).

(ii) If \( \sum_{k \geq 1} \gamma_k < +\infty \), then \( \sum_{n \geq 0} P(S_n = 0) < +\infty \).

(iii) If \((\gamma_m)\) decreases exponentially, then so does \( P(S_n = 0) \).

(iv) If \((\gamma_m)\) is summable and

\[
\alpha := \sup_i \lim_{k \to \infty} \left[ \frac{P(\tau = i)}{P(\tau = ki)} \right]^{1/k} < \frac{1}{P(\tau < +\infty)},
\]

(A.1)

then

\[
P(S_n = 0) = O(P(\tau = n)).
\]
Remark 5 We notice that, according to (5.11), \( \gamma_n \sim P(\tau = n)/P(\tau = +\infty) \). Hence, a sufficient condition for (A.1) is a similar condition for the sequence \((\gamma_n)\). In particular, statement (iv) implies the following criterion. If \((\gamma_m)\) is summable, then

\[
\sup_i \lim_{k \to \infty} \left( \frac{\gamma_i}{\gamma_{ki}} \right)^{1/k} \leq 1 \implies P(S_n = 0) = O(\gamma_n).
\]

This criterion applies, for instance, if \( \gamma_n \sim (\log n)^b n^{-a} \), for \( a > 1 \) and \( b \) arbitrary.

Sketch of the proof of (i)–(iii)

Statement (i) follows from the well known fact that the Markov chain \((S_n)_{n \in \mathbb{N}}\) is positive recurrent if and only if,

\[
\sum_{m \geq 1} \prod_{k=0}^{m} (1 - \gamma_k) < +\infty.
\]

To prove parts (ii) and (iii) we introduce the series

\[
F(s) = \sum_{n=1}^{+\infty} P(\tau = n) s^n,
\]

and

\[
G(s) = \sum_{n=0}^{+\infty} P(S_n = 0) s^n
\]

where the random variable \( \tau \) is the time of first return to zero, defined in (5.8). The probabilities \( P(\tau = n) \) were computed in (5.11) above. The relation (5.7) implies that these series are related in the form

\[
G(s) = \frac{1}{1 - F(s)},
\]

for all \( s \geq 0 \) such that \( F(s) < 1 \).

It is clear that the radius of convergence of \( F \) is at least 1. In fact,

\[
F(1) = P(\tau < +\infty).
\]

Moreover, if \( \sum_{k \geq 1} \gamma_k < +\infty \), the radius of convergence of \( F \) is

\[
\lim_{n \to \infty} [\gamma_n]^{-1/n}.
\]

This is a consequence of the fact that \( P(\tau = n)/\gamma_n^{-1} \to P(\tau = +\infty) > 0 \), as concluded from (5.11).

Statement (ii) of the proposition is a consequence of the fact that the radius of convergence of the series \( G \) is at least 1 if \( \sum_{k \geq 1} \gamma_k < +\infty \). This follows from the relation (A.5) and the fact that the right-hand side of (A.6) is strictly less than one when the chain \((S_n^{(1)})\) is transient. In fact, both sides of (ii) are necessary and sufficient conditions for the transience of the zero state.

To prove statement (iii) let us assume that \( \gamma_m \leq C\gamma^m \) for some constants \( C < +\infty \) and \( 0 < \gamma < 1 \). By (A.7), the radius of convergence of \( F \) is larger than \( \gamma^{-1} > 1 \) while, by (A.6), \( F(1) < 1 \). By continuity it follows that there exists \( s_0 > 1 \) such that \( F(s_0) = 1 \) and, hence, by (A.5), \( G(s) = +\infty \) for all \( s < s_0 \). By definition of \( G \), this implies that \( P(S_n = 0) \) decreases faster than \( \zeta^n \) for any \( \zeta \in (s_0^{-1},1) \).

Proof of (iv)
We start with the following observation. If \( i_1 + \cdots + i_k = n \), then \( \max_{1 \leq m \leq k} i_m \geq n/k \) and thus, for \( g \) increasing
\[
g(n) \leq g(k \max i_m),
\]
where \( i_{\max} = \max_{1 \leq m \leq k} i_m \). If we apply this to \( g(n) = 1/P(\tau = n) \), which is increasing by (5.11), we obtain
\[
1 \leq \frac{P(\tau = n)}{P(\tau = k i_{\max})}, \tag{A.8}
\]
We now invoke the following explicit relation between the coefficients of \( F \) and \( G \).
\[
P(S_n = 0) = \sum_{k=1}^{n} \sum_{i_1, \ldots, i_k \geq 1 \atop i_1 + \cdots + i_k = n} \prod_{m=1}^{k} P(\tau = i_m), \tag{A.9}
\]
for \( n \geq 1 \). Multiplying and dividing each factor in the rightmost product by \( P(\tau < +\infty) \), this formula can be rewritten as
\[
P(S_n = 0) = \sum_{k=1}^{n} P(\tau < +\infty)^k \sum_{i_1, \ldots, i_k \geq 1 \atop i_1 + \cdots + i_k = n} \prod_{m=1}^{k} P(\tau = i_m | \tau < +\infty). \tag{A.10}
\]
Combining this with (A.8) we obtain
\[
P(S_n = 0) \leq P(\tau = n) \sum_{k=1}^{n} P(\tau < +\infty)^k \sum_{i_1, \ldots, i_k \geq 1 \atop i_1 + \cdots + i_k = n} \prod_{m=1}^{k} \frac{P(\tau = i_m | \tau < +\infty)}{P(\tau = k i_{\max})}, \tag{A.11}
\]
If we single out the factor \( P(\tau = i_{\max} | \tau < +\infty) = P(\tau = i_{\max})/P(\tau < +\infty) \) from the rightmost product of (A.11) and use the hypothesis (A.1) we get
\[
P(S_n = 0) \leq C P(\tau = n) \sum_{k=1}^{n} \alpha^k P(\tau < +\infty)^{k-1} \sum_{i_1, \ldots, i_k \geq 1 \atop i_1 + \cdots + i_k = n} \prod_{m=1}^{k} P(\tau = i_m | \tau < +\infty), \tag{A.12}
\]
for some constant \( C > 0 \). To bound the last sum on the right-hand side we introduce a sequence of independent random variables \((\tau^{(i)})_{i \in \mathbb{N}}\) with common distribution
\[
P(\tau^{(i)} = j) = P(\tau = j | \tau < +\infty). \tag{A.13}
\]
Then
\[
\sum_{j=1}^{n-k+1} \alpha P(\sum_{i=1}^{k-1} \tau^{(i)} = n-j) \leq 1. \tag{A.14}
\]
Hence, (A.12) implies
\[
P(S_n = 0) \leq C \sum_{k=1}^{\infty} \alpha^{k-1} P(\tau < +\infty)^{k-1} P(\tau = n) \leq \text{const} P(\tau = n). \tag{A.15}
\]

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References

[1] M. BARNESLEY, S. DEMKO, J. ELTON and J. GERINOMO (1988). Invariant measures for Markov processes arising from iterated function systems with place-dependent probabilities. Ann. Inst. Henri Poincaré, Prob. Statist., 24:367–394. Erratum (1989) Ann. Inst. Henri Poincaré Probab. Statist. 25:589–590.

[2] H. BERBEE (1987). Chains with infinite connections: Uniqueness and Markov representation. Probab. Th. Rel. Fields, 76:243–253.

[3] R. BOWEN (1975). Equilibrium states and the ergodic theory of Anosov diffeomorphisms, Lecture Notes in Mathematics, volume 470. Springer.

[4] X. BRESSAUD, R. FERNANDEZ, and A. GALVES (1998). Speed of $\bar{d}$-convergence for Markov approximations of chains with complete connections. A coupling approach. To be published in Stoch. Proc. and Appl. A preliminary version can be retrieved from http://mpej.unige.ch/mparc/c/97/97-587.ps.gz.

[5] K. BURDZY and W. KENDALL (1998). Efficient Markovian couplings: examples and counterexamples. Preprint. Can be retrieved from http://www.math.washington.edu/~burdzy/Papers/list.html

[6] Z. COELHO and A. N. QUAS (1998). Criteria for $\bar{d}$-continuity. Trans. Amer. Math. Soc., 350:3257–3268.

[7] W. DOEBLIN (1938). Exposé sur la théorie des chaînes simples constantes de Markoff à un nombre fini d’états. Rev. Math. Union Interbalkanique, 2:77–105.

[8] W. DOEBLIN (1940). Remarques sur la théorie métrique des fractions continues. Composition Math., 7:353–371.

[9] W. DOEBLIN and R. FORTET (1937). Sur les chaînes à liaisons complètes. Bull. Soc. Math. France, 65:132–148.

[10] T. E. HARRIS (1955). On chains of infinite order. Pacific J. Math., 5:707–724.

[11] M. IOSIFESCU (1978). Recent advances in the metric theory of continued fractions. Trans. 8th Prague Conference on Information Theory, Statistical Decision Functions, Random Processes (Prague, 1978), Vol. A, pp. 27–40, Reidel, Dordrecht-Boston, Mass.

[12] M. IOSIFESCU (1992). A coupling method in the theory of dependence with complete connections according to Doeblin. Rev. Roum. Math. Pures et Appl., 37:59–65.

[13] M. IOSIFESCU and S. GRIGORESCU (1990). Dependence with Complete Connections and its Applications, Cambridge University Press, Cambridge.

[14] M. IOSIFESCU and R. THEODORESCU (1969). Random Processes and Learning, Springer-Verlag, Berlin.

[15] T. KAIJSER (1981). On a new contraction condition for random systems with complete connections. Rev. Roum. Math. Pures et Appl., 26:1075–1117.

[16] T. KAIJSER (1994). On a theorem of Karlin. Acta Applicandae Mathematicae, 34:51–69.

[17] S. KARLIN (1953). Some random walks arising in learning models. J. Pacific Math., 3:725–756.

[18] A. KONDAH, V. MAUME and B. SCHMITT (1996). Vitesse de convergence vers l’état d’équilibre pour des dynamiques markoviennes non höldériennes. Technical Report 88, Université de Bourgogne.

[19] S. P. LALLEY (1986). Regeneration representation for one-dimensional Gibbs states. Ann. Prob., 14:1262–1271.
[20] F. LEDRAPPIER (1974). Principe variationnel et systèmes dynamiques symboliques. Z.
Wahrscheinlichkeitstheorie verw. Gebiete, 30:185–202.

[21] T. LINDVALL (1991). W. Doeblin 1915–1940. Ann. Prob., 19:929–934.

[22] T. LINDVALL (1992). Lectures on the Coupling Method, Wiley, New York.

[23] C. LIVERANI (1995). Decay of correlations. Ann. of Math., 142(2):239–301.

[24] F. NORMAN (1972). Markov Processes and Learning Models, Academic Press, New York.

[25] O. ONICESCU and G. MIHOC (1935). Sur les chaînes statistiques. C. R. Acad. Sci. Paris,
200:5121—512.

[26] O. ONICESCU and G. MIHOC (1935a). Sur les chaînes de variables statistiques. Bull. Sci.,
Math., 59:174–192.

[27] W. PARRY and M. POLLICOTT (1990). Zeta functions and the periodic structure of hyper-
bolic dynamics. Asterisque 187-188.

[28] M. POLLICOTT (1997). Rates of mixing for potentials of summable variation. Preprint.

[29] A. N. QUAS (1996). Non-ergodicity for $c^1$ expanding maps and $g$-measures. Ergod. Th.
Dynam. Sys., 16:531–543.

[30] D. RUELLE (1978). Thermodynamic formalism. Encyclopedia of Mathematics and its appli-
cations, volume 5. Addison-Wesley, Reading, Massachusetts.

[31] P. WALTERS (1975). Ruelle’s operator theorem and $g$-measures. Trans. Amer. Math. Soc.,
214:375–387.

[32] P. WALTERS (1978). Invariant measures and equilibrium states for some mappings which
expand distances. Trans. Amer. Math. Soc., 236:121–153.

[33] L.-S. YOUNG (1997). Recurrence times and rates of mixing. Preprint.

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