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Occupancy distributions in Markov chains via Doeblin’s ergodicity coefficient

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We state and prove new properties about Doeblin’s ergodicity coefficient for finite Markov chains. We show that this coefficient satisfies a sub-multiplicative type inequality (analogous to the Markov-Dobrushin’s ergodicity coefficient), and provide a novel but elementary proof of Doeblin’s characterization of weak-ergodicity for non-homogeneous chains. Using Doeblin’s coefficient, we illustrate how to approximate a homogeneous but possibly non-stationary Markov chain of duration $n$ by independent and short-lived realizations of an auxiliary chain of duration of order $\ln(n)$. This leads to approximations of occupancy distributions in homogeneous chains, which may be particularly useful when exact calculations via one-step methods or transfer matrices are impractical, and when asymptotic approximations may not be yet reliable. Our findings may find applications to pattern problems in Markovian and non-Markovian sequences that are treatable via embedding techniques.

Keywords: compound Poisson approximation, Doeblin, Markov chain embedding technique, motif, occupancy distribution, pattern

1 Introduction.

In what follows, $S$ is a given finite set and $T \subset S$ a certain non-empty subset of states. For a fixed integer $n \geq 1$, consider a first-order homogeneous Markov chain $X = (X_t)_{0 \leq t \leq n}$ with initial distribution $\mu : S \to [0, 1]$ and probability transition matrix $p : S \times S \to [0, 1]$. We identify complex-valued functions defined over $S$ and $S \times S$ as row vectors and matrices, respectively. In particular, the distribution of $X_t$ is given by the vector $\mu p^t$.

Our object of interest is the occupancy distribution of $T$ i.e. the distribution of the random variable:

$$T_n = \sum_{i=1}^{n} [X_i \in T],$$

where $[\cdot]$ denotes the indicator function of the event within (Iverson’s bracket). Random variables of this sort are common in assessing the frequency statistics of patterns in random sequences, which typically

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model text or genomic sequences. Although various probabilistic and analytic techniques have been used for this purpose, the **Markov chain embedding technique** is among the most versatile ones. This technique seems to have originated in the works of Gerber and Li [GL81], Biggins and Cannings [BC87], and Bender and Kochman [BK93]. It consists in embedding a random sequence into the state space of a suitable finite automaton that is informative of the pattern of interest, and it has been completely systematized for **regular patterns** i.e. patterns described by a regular expression, and **Markovian models** of random sequences [NSF02, Nic03]. In addition, the technique has also shown some promise for assessing regular patterns in **non-Markovian sequences** i.e. sequences with an arbitrary correlation structure [Lla08].

All the complexity associated with determining or approximating the distribution of $T_n$ is due to the distributional dependence between the consecutive states visited by the chain. There are various ways—some more ad hoc and others more systematic—to pinpoint this distribution. For small values of $n$, exact calculations are possible via one-step methods [Dur99] or transfer matrices [FS09]. Furthermore, transfer matrices lead to Normal approximations for large values of $n$ e.g. as shown in [NSF02] for frequency statistics of regular patterns under Markovian models. On the other hand, for stationary chains, Poisson [BHJ92] and compound Poisson approximations [Erh99] have been proposed when $T$ is a rare set i.e. the stationary measure of $T$ is small. A specialized instance of these approximations is the Pólya-Aeppli distribution which occurs as the limiting distribution of frequency statistics of rare words under stationary Markov models [RS07].

Our goal is to approximate the distribution of $T_n$, when $X$ is possibly non-stationary, and in the intermediate regime where $n$ is perhaps too large for exact calculations but too small to rely on asymptotic approximations. Our approach relies on a novel use of the so called **Doeblin’s ergodicity coefficient** associated with $p$ [Doe37], which is defined as:

$$\alpha(p) \overset{def}{=} \sum_j \min_i p(i, j).$$

The above goal is far from artificial! For instance, extensive research is being performed to understand the evolution of complex but relatively short RNA sequences from simpler but functional RNA sequences [KLYK08, KLW+10]. In contexts like this, the pitfall of the Normal approximation of $T_n$ is the slow rate of convergence of order $n^{-1/2}$. On the other hand, the stationary assumption of the aforementioned Poisson approximations is unrealistic in the context of the Markov chain embedding technique, even if the background model of a genomic sequence is Markovian and stationary. For example, for regular patterns, the initial distribution of the embedded process is concentrated in a few states (of an exponentially large state space) associated with the unique initial state of the deterministic finite automaton that recognizes the pattern of interest [Lla07].

### 1.1 Ergodicity coefficients and inhomogeneous Markov chains.

This section finishes the Introduction with a brief discussion about ergodicity coefficients and the historical developments surrounding the characterization of weak-ergodicity of non-homogenous Markov chains.

In what follows we denote the set of all probability transition matrices over the state space $S$ as $\mathcal{P}$. The set of all stochastic matrices with identical rows is denoted $\mathcal{E}$; in particular, $\mathcal{E} \subset \mathcal{P}$. We refer to matrices in $\mathcal{E}$ as i.i.d. models because a homogeneous Markov chain with a probability transition matrix in this set is just a sequence of independent and identically distributed (i.i.d.) $S$-valued random variables.
Broadly speaking, an ergodicity coefficient is any continuous function \( \gamma : \mathcal{P} \rightarrow [0, 1] \). Such function is said proper when \( \gamma(p) = 1 \) if and only if \( p \in \mathcal{E} \). Clearly, Doeblin’s coefficient as defined in (2) is proper.

Other ergodicity coefficients found in the literature are:

\[
\begin{align*}
\gamma_1(p) & \overset{\text{def}}{=} \max_j \min_i p(i, j); \\
\gamma_2(p) & \overset{\text{def}}{=} \min_{i,j} \sum_s \min \{ p(i, s), p(j, s) \} = 1 - \frac{1}{2} \cdot \max_{i,j} \sum_s |p(i, s) - p(j, s)|; \\
\gamma_3(p) & \overset{\text{def}}{=} 1 - \max_s \max_{i,j} |p(i, s) - p(j, s)|.
\end{align*}
\]

Only the last two of these are proper. \( \gamma_2 \) is called the Markov-Dobrushin’s ergodicity coefficient [Mar06, Dob56a, Dob56b].

Ergodicity coefficients have been proposed for a range of purposes such as to analyze the contractive property of a stochastic matrix [Mar06] and bound its non-Perron-Frobenius eigenvalues [Sen93]. However, they have mostly been used to analyze the asymptotic behavior of non-homogeneous Markov chains [Doe37, Dob56a, Sen73b]. This entails understanding the asymptotic behavior of products of the form \( \prod_{k=m}^{n} p_k \), with \( m \leq n \), for a given sequence \( (p_k)_{k \geq 0} \subset \mathcal{P} \). Such a sequence is said to be weakly ergodic if for all \( m \geq 0 \) and \( i, j, s \in S \) the following applies:

\[
\lim_{n \to \infty} \left| \left( \prod_{k=m}^{n} p_k \right)(i,s) - \left( \prod_{k=m}^{n} p_k \right)(j,s) \right| = 0.
\]

(3)

(This condition does not imply that \( \prod_{k=m}^{n} p_k \) converges to a limit as \( n \) tends to infinity. As a counterexample, just consider the case with \( p_k = E_0 \) for \( k \) even, \( p_k = E_1 \) for \( k \) odd, for given but different matrices \( E_0, E_1 \in \mathcal{E} \).)

The following condition, known as Markov’s theorem [Sen73a], is sufficient for weak ergodicity:

\[
\sum_{k=0}^{\infty} \gamma_1(p_k) = +\infty.
\]

(4)

This condition is only sufficient in great part because \( \gamma_1 \) is not proper [Sen73b].

In more probabilistic terms, consider a first-order Markov chain \( Y = (Y_k)_{k \geq 0} \) with state space \( S \) such that \( \mathbb{P}[Y_k = s \mid Y_{k-1}, \ldots, Y_0] = p_k(Y_{k-1}, s) \), for each \( k \geq 1 \). The sequence \( (p_k)_{k \geq 0} \) is weakly ergodic if and only if any two independent realizations of \( Y \) meet infinitely often on a same state with probability one. This characterization is due to Doeblin and appeared without proof in the report [Doe37]. The way in which Doeblin proved this result is matter of speculation and it was lost with his death in World War II (see [Sen73b] for the historical developments). Furthermore, Doeblin’s report remained unnoticed for almost two decades. During this period, the following condition was proved to be both necessary and sufficient for weak ergodicity [Haj58]:

\[
\left\{ \begin{array}{l}
\text{there exists a strictly increasing sequence of positive} \\
\text{integers } (n_k)_{k \geq 0} \text{ such that: } \sum_{k=0}^{\infty} \gamma_2\left( \prod_{i=n_k}^{n_{k+1}-1} p_i \right) = +\infty.
\end{array} \right.
\]

(5)
In contrast, Doeblin’s characterization of weak ergodicity is the following [Doe37]:

\[
\{ \text{there exists a strictly increasing sequence of positive integers } (n_k)_{k \geq 0} \text{ such that: } \sum_{k=0}^{\infty} \alpha\left( \prod_{i=n_k}^{n_{k+1}-1} p_i \right) = +\infty. \}
\] (6)

Since \( \gamma_1(p) \leq \alpha(p) \leq \gamma_2(p) \), for each \( p \in \mathcal{P} \), the sufficient condition in (4) is a special instance of the conditions in (5) and (6). Though nobody knows how Doeblin proved that conditions (4) and (6) are equivalent, Seneta ventured in [Sen73b] a possible proof. His proof relies on the following two facts, valid any sequence \((p_k)_{k \geq 0} \subset \mathcal{P} :\)

(a) \( \left(1 - \gamma_3(\prod_{k=0}^{n} p_k)\right) \leq \prod_{k=0}^{n} \left(1 - \gamma_2(p_k)\right) \), for all \( n \geq 0; \)

(b) if \( \sum_{k=0}^{\infty} \alpha(p_k) = +\infty \) then \( \sum_{k=0}^{\infty} \gamma_1(p_k) = +\infty. \)

**Paper overview and organization.** Our paper is mostly about Doeblin’s ergodicity coefficient, which we encountered—by accident—while aiming at accurate but low-to-moderate complexity approximations of occupancy distributions in homogenous Markov chains. Here we mostly state and prove new properties about Doeblin’s coefficient which we would have never explored otherwise. The more detailed implications of these properties to approximate occupancy distributions will be part of a follow up publication based on the M.S. thesis [Che10].

In §2 we demonstrate new properties about Doeblin’s coefficient, which we use in §2.1 to provide a new and more elementary proof of Doeblin’s characterization of weak-ergodicity of non-homogenous Markov chains. In §3, we relate Doeblin’s coefficient to a decomposition of a chain into several independent realizations of an auxiliary chain. This decomposition has been used in the Markov chain Monte Carlo literature to sample from the stationary distribution. In §3.1 we illustrate how this decomposition allows one to approximate the distribution of a particular \( X_n \) with \( O(\ln(n)) \) rather than \( n \) matrix multiplications. In §3.2 we go a step further and illustrate how the decomposition allows us to parse (with high probability) the trajectory of a Markov chain of duration \( n \) into short-lived realizations of an auxiliary chain of durations of order \( \ln(n) \). We exploit this feature to propose new approximations for occupancy distributions based on Doeblin’s coefficient, which we compare against Normal and Poisson approximations in a numerical example.

2 A candidate for Doeblin’s missing proof.

Recall that Doeblin’s ergodicity coefficient associated with a \( p \in \mathcal{P} \) is the quantity:

\[ \alpha(p) = \sum_{j} \min_{i} p(i,j). \]

Because \( \alpha(\cdot) \) is a proper ergodicity coefficient, \( \alpha(p) \) is close to 1 when \( p \) is in the proximity of some i.i.d. model. However, since the set of i.i.d. models is closed, there should be several i.i.d. models close to \( p \). The following result identifies an affine space of i.i.d. models that are in the proximity of \( p \).
Theorem 2.1 For each \( p \in \mathcal{P} \), the following applies:

(a) If \( 0 \leq \alpha \leq \alpha(p) \) then there is \( E \in \mathcal{E} \) and \( M \in \mathcal{P} \) such that \( p = \alpha \cdot E + (1 - \alpha) \cdot M \).

(b) \( \alpha(p) = \sup \{ \alpha \in [0, 1] \mid (\exists E \in \mathcal{E})(\exists M \in \mathcal{P}); p = \alpha \cdot E + (1 - \alpha) \cdot M \} \).

(c) Assume that \( E \in \mathcal{E} \) and \( M \in \mathcal{P} \) are such that \( p = \alpha(p) \cdot E + (1 - \alpha(p)) \cdot M \).

If \( \alpha(p) < 1 \) then \( \alpha(M) = 0 \) i.e. \( M \) has a zero in each column.

If \( \alpha(p) > 0 \) then \( E(i, j) = \frac{1}{\alpha(p)} \cdot \min_s p(s, j) \).

Proof: Define \( \beta = \alpha(p) \). We first show part (a) in the theorem, for which we may assume without loss of generality that \( \beta > 0 \). In this case, all reduces to prove that there is \( E \in \mathcal{E} \) and \( M \in \mathcal{P} \) such that

\[
 p = \beta \cdot E + (1 - \beta) \cdot M. 
\]  

(7)

Indeed, if \( 0 \leq \alpha \leq \beta \) then the above implies that \( p = \alpha E + (\beta - \alpha) E + (1 - \beta) M = \alpha E + (1 - \alpha) Q \), for some \( Q \in \mathcal{P} \), because the matrix \((\beta - \alpha) E + (1 - \beta) M\) has nonnegative entries and the sum of the entries in each of its rows is \((1 - \alpha)\). To prove the above identity, consider the matrix \( E \in \mathcal{E} \) with entries \( E(i, j) = \min_s p(s, j)/\beta \). Since \( \beta E(i, j) \leq p(i, j) \), the matrix \((p - \beta E)\) has nonnegative entries and row sums equal to \((1 - \beta)\). In particular, if \( \beta = 1 \) then \( p = E \) and the above identity holds with any choice of \( M \). Otherwise, it suffices to select \( M = (p - \beta E)/(1 - \beta) \). This shows (7) and completes the proof of part (a).

To show part (b), notice that if \( 0 \leq \alpha \leq 1 \) is such that \( p = \alpha E + (1 - \alpha) M \), with \( E \in \mathcal{E} \) and \( M \in \mathcal{M} \), then \( \beta = \alpha(p) \geq \alpha \cdot \alpha(E) = \alpha \). Part (b) is now a direct consequence of part (a).

Finally we show part (c). Thus assume that \( E \in \mathcal{E} \) and \( M \in \mathcal{P} \) are such that \( p = \beta E + (1 - \beta) M \). If \( \beta = 1 \) then \( p = E \) and the identity \( E(i, j) = \min_s p(s, j)/\beta \) is trivial. On the other hand, if \( \beta = 0 \) then \( p \) must have a zero in each column and \( M = p \). Without loss of generality we may therefore assume that \( 0 < \beta < 1 \). We first show that \( \alpha(M) = 0 \). Set \( \alpha' = \alpha(M) \). Due to part (a), there exists \( E' \in \mathcal{E} \) and \( M' \in \mathcal{P} \) such that \( p = \beta E + (1 - \beta) \alpha' E' + (1 - \beta)(1 - \alpha') M' \). Hence \( \beta = \alpha(p) \geq (\beta + (1 - \beta) \alpha') \) and as a result \( \alpha' = 0 \). To complete the proof of the theorem, fix \( j \) and notice that \( \beta E(i, j) = p(i, j) - (1 - \beta) M(i, j) \). In particular, since \( M \) has a zero in each column, there is \( s \) such that \( \beta E(s, j) = p(s, j) \). Finally, since \( \beta E(i, j) \leq p(i, j) \), we conclude that \( \beta E(i, j) = \min_s p(s, j) \). This completes the proof. \( \square \)

Due to part (a) in the previous theorem, for all \( p \in \mathcal{P} \), there is \( E \in \mathcal{E} \) and \( M \in \mathcal{P} \) such that:

\[
(p - E) = (1 - \alpha(p)) \cdot (M - E).
\]

As a result, the smaller \((1 - \alpha(p))\), the closer \( p \) is to an i.i.d. model. According to the following corollary, \((1 - \alpha(\cdot))\) is a sub-multiplicative function; in particular, when one multiplies two or more stochastic matrices, one can only get “closer” to the set of i.i.d. models. This is the key ingredient for our proof of Doeblin’s characterization of weak ergodicity.

To the best of our knowledge, the sub-multiplicative property has been shown in the literature only for the Markov-Dobrushin ergodicity coefficient [Dob56a, Dob56b, Ios72, Paz70, Gri75].

Corollary 2.2 \((1 - \alpha(pq)) \leq (1 - \alpha(p)) \cdot (1 - \alpha(q))\), for all \( p, q \in \mathcal{P} \).
**Proof:** Define $\alpha_1 = \alpha(p)$ and $\alpha_2 = \alpha(q)$. Due to part (a) in Theorem 2.1, there are matrices $E_1, E_2 \in \mathcal{E}$ and $M_1, M_2 \in \mathcal{P}$ such that $p = \alpha_1 E_1 + (1 - \alpha_1) M_1$ and $q = \alpha_2 E_2 + (1 - \alpha_2) M_2$. In particular, since $p E_2 = E_2$, we see that $pq = \alpha_2 E_2 + \alpha_1 (1 - \alpha_2) E_1 M_2 + (1 - \alpha_1) (1 - \alpha_2) M_1 M_2$. But notice that $E_1 M_2 \in \mathcal{E}$. Consequently, the rows of the matrix $\alpha_2 E_2 + \alpha_1 (1 - \alpha_2) E_1 M_2$ are identical, with common row sum $(\alpha_1 + \alpha_2 - 1)$. As a result, there is $E_3 \in \mathcal{E}$ such that

$$pq = (\alpha_1 + \alpha_2 - \alpha_1 \alpha_2) \cdot E_3 + (1 - \alpha_1)(1 - \alpha_2) \cdot M_1 M_2;$$

$$= (1 - (1 - \alpha_1)(1 - \alpha_2)) \cdot E_3 + (1 - \alpha_1)(1 - \alpha_2) \cdot M_1 M_2.$$

Finally, due to part (b) in Theorem 2.1, it follows from the above that

$$\alpha(pq) \geq 1 - (1 - \alpha_1)(1 - \alpha_2),$$

which proves the corollary. \( \square \)

## 2.1 A first principles proof of Doeblin’s characterization of weak ergodicity.

As we mentioned earlier, Doeblin’s proof that (3) and (6) are equivalent is matter of speculation. Though it is possible to prove this equivalence using Theorem 1 in [Sen73] and Corollary 2.2 in the previous section, we venture an alternative and more elementary proof of this fact.

Fix an integer $m \geq 0$ and let $\alpha_n$ denote the Doeblin’s ergodicity coefficient of $\prod_{k=m}^{n} p_k$. Due to parts (a) and (c) in Theorem 2.1, there are matrices $E_n \in \mathcal{E}$ and $M_n \in \mathcal{P}$ such that:

$$\prod_{k=m}^{n} p_k = \alpha_n \cdot E_n + (1 - \alpha_n) \cdot M_n,$$

with $\alpha(M_n) = 0$.

In particular, for all $i, j, s \in S$ the following holds:

$$(\prod_{k=m}^{n} p_k)(i, s) - (\prod_{k=m}^{n} p_k)(j, s) = (1 - \alpha_n) \cdot (M_n(i, s) - M_n(j, s)). \quad (8)$$

Assume first that condition (6) holds. Consider the sets of non-negative integers:

$$I_n = \{k \mid \exists j \text{ such that } m \leq n_j \leq k \leq n_{j+1} \leq n\};$$

$$J_n = \{j \mid \{n_j, \ldots, n_{j+1}\} \subset I_n\}.$$ 

Notice that $J_n$ is an interval of integers. Furthermore, there are stochastic matrices $L$ and $R_n$ such that $\prod_{k=m}^{n} p_k = L \cdot (\prod_{k \in I_n} p_k) \cdot R_n$. In particular, due to Corollary 2.2, we find that

$$(1 - \alpha_n) \leq (1 - \alpha (\prod_{k \in I_n} p_k)) \leq \prod_{j \in J_n} (1 - \alpha (\prod_{k=n_j}^{n_{j+1} - 1} p_k)).$$

Since $(1 - x) \leq \exp(-x)$, for $0 \leq x \leq 1$, the condition in (6) implies that $\lim_{n \to \infty} \alpha_n = 1$. Back in (8), since each $M_n$ is a stochastic matrix, we conclude that

$$\lim_{n \to \infty} (\prod_{k=m}^{n} p_k)(i, s) - (\prod_{k=m}^{n} p_k)(j, s) = 0.$$
This shows that condition (3) is also satisfied i.e. \((p_k)_{k \geq 0}\) is weakly ergodic.

Conversely, assume that condition (3) holds. To show that condition (6) also applies, we first prove that \((\alpha_n)_{n \geq 0}\) has a subsequence that converges to one. We show this by contradiction. In particular, due to the identity in (8), it applies that

\[
\lim_{n \to \infty} \left( M_n(i, s) - M_n(j, s) \right) = 0,
\]

for all \(i, j, s \in S\). Fix \(s_1 \in S\). Since each \(M_n\) has at least one zero in the column associated with \(s_1\) then there is \(j_1 \in S\) and a subsequence \((n'_k)_{k \geq 0}\) such that \(M_{n'_k}(j_1, s_1) = 0\) for all \(k \geq 0\). Therefore, \(M_{n'_k}(i, s_1) \to 0\) as \(k \to \infty\), for all \(i \in S\). Now fix \(s_2 \in S \setminus \{s_1\}\). Since each \(M_{n'_k}\) has at least one zero in the column associated with \(s_2\) then there is a subsequence \((n''_k)_{k \geq 0}\) of \((n'_k)_{k \geq 0}\) such that \(M_{n''_k}(i, s_1) \to 0\) and \(M_{n''_k}(i, s_2) \to 0\) as \(k \to \infty\), for all \(i \in S\). Since \(S\) is finite, a straightforward inductive argument shows that there is a subsequence \((n_k)_{k \geq 0}\) such that

\[
\lim_{k \to \infty} M_{n_k}(i, s) = 0,
\]

for all \(i, s \in S\). However, the above is not possible because each \(M_{n_k}\) is a stochastic matrix. As a result, \((\alpha_n)_{n \geq 0}\) must have a subsequence that converges to one.

The previous argument shows that if \((p_k)_{k \geq 0}\) is weakly ergodic then for all \(m \geq 0\) there is \(n \geq m\) such that e.g. \(\alpha(\prod_{k=m}^{n} p_k) \geq 1/2\). From this, condition (6) is immediate and we have proved that conditions (3) and (6) are equivalent.

3 Occupancy distributions of homogeneous chains.

In this section, we retake our initial motivation of approximating occupancy distributions of homogeneous Markov chains. As a starting point, suppose that:

\[
p = \alpha \cdot E + (1 - \alpha) \cdot M,
\]

for certain \(0 \leq \alpha \leq 1\), \(E \in \mathcal{E}\) and \(M \in \mathcal{P}\), and denote as \(e\) any of the rows of \(E\) (see Theorem 2.1). The following stochastic equivalent of \(X = (X_i)_{i \geq 0}\) has been exploited in the Markov chain Monte Carlo literature to sample—perfectly—from the stationary distribution, without having to compute it beforehand [MG98, Møl99, CT01].

**Lemma 3.1** Assume that condition (9) is satisfied. Imagine a coin that shows \(E\) with probability \(\alpha\) and \(M\) with probability \((1 - \alpha)\) when tossed. The stochastic sequence \(Y = (Y_i)_{i \geq 0}\) defined as follows:

(i) \(Y_0\) has distribution \(\mu\), and

(ii) for each \(i \geq 0\), the distribution of \(Y_{i+1}\) conditioned on \((Y_0, \ldots, Y_i)\) is given by the following procedure: toss the coin, and if the \(E\)-side comes up then draw \(Y_{i+1}\) using the distribution \(e(\cdot)\), else draw \(Y_{i+1}\) using the distribution \(M(Y_i, \cdot)\).

has the same distribution as \(X\).
**Proof:** Due to the definition of the $Y$ process, for each $i \geq 0$ and $s_0, \ldots, s_i+1 \in S$, the following applies:

$$
P(Y_{i+1} = s_{i+1} \mid Y_0 = s_0, \ldots, Y_i = s_i) = \alpha \cdot e(s_{i+1}) + (1 - \alpha) \cdot M(s_i, s_{i+1}),$$

$$= \alpha \cdot E(s_i, s_{i+1}) + (1 - \alpha) \cdot M(s_i, s_{i+1}) = p(s_i, s_{i+1}).$$

In particular, $Y$ is a first-order homogeneous Markov chain with initial distribution $\mu$ and probability transition matrix $p$. Hence $X$ and $Y$ have the same distribution. \hfill $\square$

Rather than simulating perfectly from the stationary distribution, we illustrate how to use Lemma [3.1] to approximate occupancy distributions in homogeneous Markov chains. The key idea is to use Doeblin’s ergodicity coefficient to break—at random times—the dependence of the chain. As an intermediate step, we consider a simpler problem: how to approximate the distribution of a particular $X_n$. Of course, if the chain is aperiodic and irreducible then, for large $n$, the distribution of $X_n$ will be well-approximated by the stationary distribution. However, as pointed out in [Doe38], our goal is to provide approximations in the intermediate regime where asymptotic approximations may not yet be reliable.

For what follows, recall that the total variation distance between two probability distributions $u(\cdot)$ and $v(\cdot)$ supported over $\mathbb{N} = \{0, 1, \ldots\}$ is defined as:

$$\|u - v\| \overset{def}{=} \sup_{A \subset \mathbb{N}} |u(A) - v(A)| = \frac{1}{2} \sum_{i \in \mathbb{N}} |u(\{i\}) - v(\{i\})|.$$ 

Accordingly, the total variation distance between two $\mathbb{N}$-valued random variables $U$ and $V$ is defined as the total variation distance of their distributions and it is denoted $\|U - V\|$.

### 3.1 A simpler problem.

It is well-known that, in finite state spaces, if $p$ is irreducible and aperiodic then there is a unique stationary distribution i.e. a unique probability distribution $\pi$ such that $\pi p = \pi$. In this case, there are constants $c_0, c_1 > 0$ which depend on $p$ but not on $\mu$, such that:

$$\|X_n - \pi\| \leq c_0 e^{-c_1 n}, \text{ for all } n \geq 0. \quad (10)$$

The proof of this fact can be found in most books on Markov chain theory, with approaches ranging from the well-known Perron-Frobenius Theorem [Per07, Fro12] to Operators and Ergodic theory [Lin71]. Various proofs are based on the coupling method, generally attributed to Doeblin, who proved the above in [Doe38]. The coupling technique has been since then refined to address Markov chains with countable state spaces [Pit74, Gri75, Tho90] and a purely probabilistic proof of the Renewal theorem [Lin77].

Using Lemma [3.1] we may alternatively explain (10) as follows. If $\alpha(p) > 0$ then there is a distribution $\epsilon$ such that, regardless of the state where the chain is located, the next state will be picked up from this distribution with probability $\alpha(p)$. Each time this distribution is used, any information about the states previously visited by the chain is lost. This distribution acts therefore as a “memory-breaker”. When $n$ is large, and even if $\alpha(p)$ is small, it is unlikely that no memory-breaker occurred between $X_0$ and $X_n$. Since all transitions after the last memory-breaker where controlled by $M$, the distribution of $X_n$ should be well-approximated by a mixture of distributions of the form $\epsilon M^t$. This intuition is made precise on the following result. Due to part (b) in Theorem [2.1] notice that the optimal choice for $\alpha$ is $\alpha(p)$. 


Corollary 3.2. [Che10] Assume that condition (9) is satisfied. Let \( \alpha \geq 0 \) and consider \( S \)-valued random variables \( Z_0, \ldots, Z_m \) such that \( Z_t \) has distribution \( e^{MT} \). If \( I \) is a random index independent of \( (Z_1, \ldots, Z_m) \) such that \( P[I = t] = \alpha (1 - \alpha)^{t-1} (1 - (1 - \alpha)^m / m + 1) \), for \( 0 \leq t \leq m \); in particular, \( P[I \in \{0, \ldots, m\}] = 1 \), then
\[
\|X_n - Z_I\| \leq (1 - \alpha)^{m+1}, \text{ for all } n > m.
\]

To fix ideas, consider the probability transition matrix:
\[
p = \begin{bmatrix}
\frac{3}{10} & 0 & \frac{7}{10} \\
0 & \frac{9}{10} & \frac{1}{10} \\
\frac{4}{5} & \frac{1}{5} & 0
\end{bmatrix}.
\]

Since \( \alpha(p) = 0 \), the inequalities of the corollary are trivial. However, observe that:
\[
p^2 = \begin{bmatrix}
\frac{13}{20} & \frac{7}{20} & \frac{21}{100} \\
\frac{2}{25} & \frac{83}{100} & \frac{9}{100} \\
\frac{6}{25} & \frac{9}{25} & \frac{29}{100}
\end{bmatrix} = \frac{31}{100} \cdot E_2 + \frac{69}{100} \cdot M_2,
\]

with
\[
E_2 := \begin{bmatrix}
\frac{8}{31} & \frac{14}{31} & \frac{9}{31} \\
\frac{8}{31} & \frac{14}{31} & \frac{9}{31} \\
\frac{8}{31} & \frac{14}{31} & \frac{9}{31}
\end{bmatrix} \quad \text{and} \quad M_2 := \begin{bmatrix}
\frac{19}{23} & 0 & \frac{4}{23} \\
0 & 1 & 0 \\
\frac{16}{69} & \frac{4}{69} & \frac{49}{69}
\end{bmatrix}.
\]

Notice that \( \alpha_2 := \alpha(p^2) = 31/100 \); in particular, the above decomposition of \( p^2 \) is a direct consequence of part (c) in Theorem 2.1. Define \( e_2 \) as the first row of \( E_2 \). Imagine you would like to approximate the distribution of some \( X_n \), with as few matrix multiplications as possible, and within a 5% accuracy in total variation distance. Define \( \epsilon := 0.05 \). Due to the Corollary 3.2, this is possible for any even number \( n \geq 18 \), by considering a mixture of the distributions \( e_2 M_2^t \), with \( t = 0, \ldots, 8 \). On the other hand, because Markovian kernels are contractive in total variation distance [Mar06], this is also possible for any odd number \( n \geq 18 \) by considering a mixture of the distributions \( e_2 M_2^t p \), with \( t = 0, \ldots, 8 \). Either mixture can be computed in at most 10 matrix multiplications, however, as it follows from Table 1, this number can be optimized to a minimum of 7 matrix multiplications by considering larger powers of \( p \).

According to the data displayed in Table 1, it is possible to approximate within \( \epsilon \)-units the distribution of each \( X_n \), with \( n \geq 16 \), using a mixture of three distributions associated with Doeblin’s ergodicity coefficient of \( p^4 \). This mixture is given by:
\[
\sum_{t=0}^{3} \frac{(1 - \alpha_4)^t - (1 - \alpha_4)^{t+1}}{1 - (1 - \alpha_4)^4} \cdot e_4 \cdot M_4^t \cdot p^{n \pmod 4},
\]
which can be computed using no more than 7 matrix multiplications. In retrospect, this is far from obvious.

For instance, one finds that:
\[
p^7 = \begin{bmatrix}
0.3444507000 & 0.3440640000 & 0.3114853000 \\
0.1966080000 & 0.6114381000 & 0.1919539000 \\
0.3559832000 & 0.3839078000 & 0.2601090000
\end{bmatrix}.
\]
Since $\max_{i,j} \|p^7(i,\cdot) - p^7(j,\cdot)\| \geq 0.26$, the chain is still far from the stationary distribution even after 7 transitions. Indeed, $\max_{i} \|p^n(i,\cdot) - \pi(\cdot)\| \geq 0.13$, exceeding the total variation distance between any $X_n$, with $n \geq 16$, and the distribution in (12).

<table>
<thead>
<tr>
<th>$k$</th>
<th>$\alpha_k$</th>
<th>$m_k$</th>
<th>$n_k$</th>
<th>$c_k$</th>
</tr>
</thead>
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<td>0</td>
<td>$\infty$</td>
<td>$\infty$</td>
<td>$\infty$</td>
</tr>
<tr>
<td>2</td>
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<td>8</td>
<td>18</td>
<td>10</td>
</tr>
<tr>
<td>3</td>
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<td>5</td>
<td>18</td>
<td>8</td>
</tr>
<tr>
<td>4</td>
<td>0.5287</td>
<td>3</td>
<td>16</td>
<td>7</td>
</tr>
<tr>
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<td>3</td>
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<td>8</td>
</tr>
<tr>
<td>6</td>
<td>0.679891</td>
<td>2</td>
<td>18</td>
<td>8</td>
</tr>
<tr>
<td>7</td>
<td>0.7326259</td>
<td>2</td>
<td>21</td>
<td>9</td>
</tr>
</tbody>
</table>

**Tab. 1**: Parameters associated with powers of the probability transition matrix in (11). Here $\alpha_k = \alpha(p^k)$ and $p^k = \alpha_k E_k + (1 - \alpha_k) M_k$, with $E_k \in E$ and $M_k \in P$. In addition, $m_k := \ln(\epsilon)/\ln(1 - \alpha_k) - 1$, with $\epsilon := 0.05$, and $n_k := k(m_k + 1)$. Due to Corollary 3.2, if $e_k$ denotes any of the rows of $E_k$ then for each $n \geq n_k$ there exists a mixture of the distributions $e_k M^t_k p^{\text{lower}(k)}; t = 0, \ldots, m_k$, which approximates the distribution of $X_n$ within $\epsilon$-units in total variation distance. This mixture can be computed with at most $c_k := (k + m_k)$ matrix multiplications.

### 3.2 Approximation of occupancy distributions.

Assume that condition (9) is satisfied. Following the notation of Lemma 3.1, the occupancy distribution of a set $T \subset S$ is the distribution of the random variable:

$$T_n = \sum_{t=1}^{n} [Y_t \in T].$$

The moment generating function (m.g.f.) of $T_n$ is given by $\mu \cdot \{p(z)\}^n \cdot 1$, where $p(z)$ is the matrix with polynomial entries given by $p(z)(i,j) = p(i,j) \cdot z^{\#(j \in T)}$, and 1 is a column-vector of ones. $p(z)$ is called a transfer matrix [FS09], and the computation of the exact distribution of $T_n$ is expensive unless $n$ is relatively small. In what follows, we extend the argument of the previous section to approximate this distribution.

Notice that the random variables $[Y_i \in T]$ and $[Y_j \in T]$, with $i < j$, are independent when at least one of the random variables $Y_{i+1}, \ldots, Y_j$ is drawn from the the memory-breaker distribution $e$. In particular, the times at which the $E$-side of the coin appears cut the trajectory $Y_0, \ldots, Y_n$ into independent “pieces”. The number of such pieces is random, and consecutive transitions in each piece are governed by the matrix $M$. Furthermore, the initial distribution of each piece is $e$, except for the first piece which has initial distribution $\mu$. The expected number of memory-breakers between the first and last transition is $\alpha n$; and the average separation between consecutive memory-breakers is $1/\alpha$, regardless of $n$. As a result, a mixture of $e(z) \cdot \{M(z)\}^m \cdot 1$, with $m$ an integer in a neighborhood of $1/\alpha$, should lead to a decent approximation of the m.g.f. of the occupancy distribution of $T$ in each piece other than the first. For the first piece, the m.g.f.’s to consider are of the form $\mu \cdot \{M(z)\}^m \cdot 1$. Since the behavior of the Markov chain is independent from one piece to another, an approximation for the m.g.f. of $T_n$ should follow. More importantly for computations, a power of order $o(n)$ of the transfer matrix $M(z)$ should suffice for a decent approximation of the distribution of $T_n$. The weakest point of this heuristic is the probable
occurrence of longer than expected pieces at already intermediate values of $n$. This motivates us to look at the random variable $L_n$, defined as the length of the longest piece. (In probabilistic terms, $L_n$ is the
length of the largest run of $M$'s in $n$-tosses of the coin from Lemma 3.1.) The asymptotic distribution of
this random variable is well understood, both via combinatorial and probabilistic methods [Fel68, FS09, AGG90]. Since the distribution of $L_n$ concentrates around $-\ln(\alpha n)/\ln(1 - \alpha)$ as $n$ increases, selecting
$m = -c \ln(\alpha n)/\ln(1 - \alpha)$, for $c > 1$, gives $P[L_n \leq m] = 1 + O(n^{-1})$. An explicit upper-bound for the
error in total variation distance follows now from the next result. We notice that the m.g.f. of the random
variable $W_I$ in the corollary can be computed explicitly via a symbolic specification [FS09].

**Corollary 3.3** [Che10] Assume that condition (9) is satisfied. Fix $m \geq 0$ and define $\ell_n = P[L_n \leq m]$. Let $I = (I_0, \ldots, I_K)$ be a random composition of length $n$ such that $P[I = (i_0, i_1, \ldots, i_k)] = \alpha^k (1 - \alpha)^{n-k}/\ell_n$, for all $k \geq 0$, $0 \leq i_0 \leq m$, $1 \leq i_l \leq (m+1)$, for $l \geq 1$, and such that $\sum_{l=0}^{K} i_l = n$. In addition, consider independent random variables $U(l)$, $V(i,l)$ which are independent of $I$ and such that $U(l)$ has m.g.f. $\mu \cdot \{M(z)\}^l \cdot 1$ and $V(i,l)$ has m.g.f. $\epsilon(z) \cdot \{M(z)\}^{l-1} \cdot 1$. If one defines $W_I := U(I_0) + \sum_{l=1}^{K} V(I_l, I_l)$ then

$$\|T_n - W_I\| \leq (1 - \ell_n).$$

As a numerical example, we select a stationary and homogenous Markov chain from [Erh99] with state space $S = \{1, \ldots, 8\}$ and probability transition matrix

$$p(i,j) = \begin{cases} \frac{1 - \beta q(i,T)}{1 - q(i,T)}, & i \in T^c, j \in T^c; \\ \beta q(i,j), & i \in T^c, j \in T; \\ q(i,j), & i \in T, j \in S; \end{cases}$$

where

$$q = \begin{pmatrix} 0.334 & 0.215 & 0.173 & 0.119 & 0.065 & 0.086 & 0.003 & 0.005 \\ 0.289 & 0.133 & 0.211 & 0.133 & 0.067 & 0.156 & 0.007 & 0.004 \\ 0.356 & 0.184 & 0.075 & 0.043 & 0.151 & 0.183 & 0.002 & 0.006 \\ 0.41 & 0.162 & 0.108 & 0.075 & 0.14 & 0.097 & 0.005 & 0.003 \\ 0.316 & 0.239 & 0.044 & 0.218 & 0.076 & 0.098 & 0.004 & 0.005 \\ 0.44 & 0.176 & 0.044 & 0.242 & 0.088 & 0.005 & 0.006 & 0.005 \\ 0.18 & 0.06 & 0.19 & 0.09 & 0.13 & 0.1 & 0.13 & 0.12 \\ 0.2 & 0.16 & 0.07 & 0.1 & 0.14 & 0.1 & 0.09 & 0.14 \end{pmatrix}.$$
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\begin{tabular}{cccc}
\hline
\(n\) & \(\beta\) & Normal & Poisson & Compound Poisson & Approximation in [13] \\
\hline
10 & 1 & 1.7E-2 & 1.4E-2 & 3.2E-3 & 3.8E-4 \\
10 & 0.5 & 1.7E-2 & 7.0E-3 & 1.2E-3 & 1.5E-4 \\
10 & 0.25 & 1.3E-2 & 3.6E-3 & 4.9E-4 & 6.9E-5 \\
10 & 0.1 & 5.3E-3 & 1.4E-3 & 1.7E-4 & 2.7E-5 \\
10 & 0.01 & 5.3E-4 & 1.4E-4 & 1.5E-5 & 2.6E-6 \\
10 & 0.001 & 5.3E-5 & 1.4E-5 & 1.5E-6 & 2.6E-7 \\
100 & 1 & 0.23 & 6.9E-2 & 9.7E-3 & 2.3E-4 \\
100 & 0.5 & 0.22 & 5.2E-2 & 3.5E-3 & 1.6E-4 \\
100 & 0.25 & 0.14 & 3.2E-2 & 1.3E-3 & 7.5E-5 \\
100 & 0.1 & 2.0E-2 & 1.5E-2 & 3.1E-4 & 3.1E-5 \\
100 & 0.01 & 5.2E-3 & 1.6E-3 & 1.6E-5 & 3.3E-6 \\
100 & 0.001 & 5.3E-4 & 1.6E-4 & 1.5E-6 & 3.3E-7 \\
1000 & 1 & 6.9E-2 & 7.0E-2 & 9.4E-3 & 2.1E-5 \\
1000 & 0.5 & 9.0E-2 & 7.3E-2 & 4.9E-3 & 1.4E-5 \\
1000 & 0.25 & 0.14 & 7.8E-2 & 2.7E-3 & 8.2E-6 \\
1000 & 0.1 & 0.23 & 6.8E-2 & 9.6E-4 & 1.1E-5 \\
1000 & 0.01 & 2.0E-2 & 1.5E-2 & 2.7E-5 & 1.8E-6 \\
1000 & 0.001 & 5.2E-3 & 1.5E-3 & 1.7E-6 & 2.0E-7 \\
\hline
\end{tabular}

\textbf{Tab. 2:} Total variation distance for approximations to the occupancy distribution of the set \(T = \{8\}\) for the stationary chain described by [14]. The compound Poisson approximation, given by [Erh99], is a Pólya-Aeppli distribution.

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\section*{References}


Occupancy distributions via Doeblin’s ergodicity coefficient


