



Speed of \bar{d} -convergence for Markov approximations of chains with complete connections. A coupling approach[☆]

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Abstract

We compute the speed of convergence of the canonical Markov approximation of a chain with complete connections with summable decays. We show that in the \bar{d} -topology the approximation converges at least at a rate proportional to these decays. This is proven by explicitly constructing a coupling between the chain and each range- k approximation. © 1999 Elsevier Science B.V. All rights reserved.

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1. Introduction

The main result of this paper is an estimation of the speed of convergence – in the \bar{d} -distance – of the canonical Markov approximation of chains with complete connections. If the continuity rates of the chain are summable, we show that the speed of convergence is at worst proportional to these rates.

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Approximations schemes are essential for understanding and handling non-Markovian processes. The speed of convergence is perhaps the most important characterization of an approximation scheme. On the one hand, it may carry information about regularity properties of the target process. On the other hand, it can be used as a tool to design efficient numerical approaches, and to establish tests to determine whether a given process is of some particular type. These facts could be all the more relevant in relation with some strongly non-Markovian processes and fields of recent interest (see e.g. van Enter et al., 1993). Nevertheless, published results on non-Markovian random processes deal only with the issue of *existence* of Markov approximations, and properties inherited from this fact. There appears to be no result so far on speed of convergence.

The existence results apply to stationary processes that either

- (a) are the \bar{d} -limit of k -step Markov processes, or
- (b) have a continuous dependence on past history;

where \bar{d} is the distance introduced by Ornstein (see Definition 3 below and Ornstein (1974) for more details).

Stationary processes of type (a) inherit the property of being Bernoulli if the approximating Markov chains are aperiodic (Friedman and Ornstein, 1970). The use of the distance \bar{d} is definitory. Indeed, *every* process can be approximated in the *vague* topology by the so-called *canonical k -step Markov approximations*, defined so to have the same transitions from k to $k+1$ states as the original process (Definition 2 below). This fact, however, is of little use, because weak limits do not convey information on long-range properties. A more revealing issue is whether the canonically defined Markov chains provide also an approximation scheme in the finer \bar{d} topology. The class of processes for which this is true has been completely characterized by Rudolph and Schwarz (1977). In particular, totally ergodic processes have this property if and only if they are Bernoulli (Friedman and Ornstein, 1970).

Stationary processes of type (b) have been studied under the stronger hypothesis of log-continuity. Following Lalley (1986), we shall call them *chains with complete connections* (Lalley's definition differs from the one introduced by Doeblin and Fortet (1937)). Each process with exponential rates of (log-)continuity is in correspondence with the unique Gibbs measure of a one-dimensional system with an exponentially decaying interaction. If the continuity rates are summable, the process is weak Bernoulli (Ledrappier, 1974). This implies, by Ornstein theorem (Ornstein, 1974), that the process is the \bar{d} -limit of its canonical k -step Markov approximations. Curiously, this indirect argument appears to be the only published proof of such \bar{d} -convergence. In contrast, our construction below yields an explicit and direct proof.

We mention two further developments. Lalley (1986) has proposed a regenerative representation of chains with complete connections, in terms of what he calls *list processes*. These are processes which at some random times “forget the past” and “begin from scratch”. The distribution of these random times depends on the continuity rates of the initial process. It has a finite exponential moment if the rates are exponential, and only moments up to some finite order if the continuity rates decay as a power-law. On the other hand, Ornstein and Weiss (1990) have constructed a remarkable “guessing scheme” for \bar{d} -limits of aperiodic Markov processes, based on observed data.

Nevertheless, these approaches do not shed light on “how well” the chains can be approximated by Markov processes.

In this paper we analyze precisely this issue for the chains with complete connections and the less sophisticated of the approximation schemes: the canonical k -step Markov. Our results show that the continuity rates of the chain directly determine – in the summable case – an upper bound on the speed of convergence of the approximation.

Our method is constructive and straightforward. We exhibit explicit couplings between the original chain and each of its k -step approximations. The couplings are such that: (i) if the two component processes have been equal for a certain number of steps, there is a large probability that they will remain so in the next step [formula (17)], and (ii) if the components fail to be equal at some step there is a nonzero probability that they will become equal at the next one [formula (18)]. As a consequence, the coupled processes tend to coincide most of the time, and separations do not last too long [formula (22)]. This yields a small \bar{d} -distance between the original process and its k -step approximations.

We conjecture that our result is optimal in the following sense. Given a decreasing and summable sequence, there is a chain with complete connections for which this sequence gives the continuity rates and, at the same time, the exact rates of convergence of the canonical k -step Markov approximations.

Analogous questions can be posed for long range Gibbsian fields. We expect the answers to be similar to those presented here, at least at low temperature. In fact, we expect the corresponding proof to follow from a construction similar to the coupling used here.

The coupling concept was introduced by Doeblin in 1938 in a hardly known paper published at the *Revue Mathématique de l'Union Interbalkanique*. To study the convergence to equilibrium of a Markov chain, Doeblin let two independent trajectories of the process 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. For a description of Doeblin's contributions to probability theory we refer the reader to Lindvall (1991). The idea was only exploited much later, in the sixties, in papers by Athreya, Ney, Harris, Spitzer and Toom among others. Liggett (1993) reviews the use of the coupling technique for interacting Markov systems. For a nice presentation of the idea of coupling related to Chen-Stein method, we refer the reader to Barbour et al. (1992). The basic idea of our coupling can be traced back to Dobrushin (1956), even when there is no coupling in his paper. Other source of inspiration is Harris' graphical method (Harris, 1978). For a pedestrian derivation of Dobrushin's ergodic coefficient using coupling we refer the reader to Ferrari and Galves (1997). A coupling approach related to ours has been used by Marton (1996).

A problem related to the discussion of the present paper is the determination of the relaxation rate of the chain. In a recent paper, Kondah et al. (1996) have estimated this rate for one-dimensional Gibbsian systems, for non-Hölder potentials, using the technique of projective metrics. In a forthcoming paper (Bressaud et al., 1999) we shall show that similar results can be obtained using our coupling approach.

The paper is organized as follows. The main result and relevant definitions are stated in Section 2 while the proof is developed in Section 3.

2. Definitions and main result

Let $X = (X_n)_{n \in \mathbf{Z}}$ be a stationary stochastic process, defined on a probability space $(\Omega, \mathcal{F}, \mathbf{P})$, taking values in a finite set A (the “alphabet”).

Definition 1. We shall say that the process $(X_n)_{n \in \mathbf{Z}}$ is a *chain with complete connections* if it satisfies the following three properties:

- for all $a_1, \dots, a_n \in A$,

$$\mathbf{P}(X_1 = a_1, \dots, X_n = a_n) > 0 \quad (1)$$

- the limit

$$\lim_{m \rightarrow \infty} \mathbf{P}(X_0 = a_0 | X_j = a_j, -m \leq j \leq -1) = \mathbf{P}(X_0 = a_0 | X_j = a_j, j \leq -1) \quad (2)$$

exists for all $a_j, j \leq -1$,

- there is a sequence $(\gamma_m)_{m \geq 1}$ with $\lim_{m \rightarrow \infty} \gamma_m = 0$, such that, for all $a_j, b_j \in A, j \leq -1$ with $a_j = b_j$ for $-1 \geq j \geq -m$,

$$\left| \frac{\mathbf{P}(X_0 = a_0 | X_j = a_j, j \leq -1)}{\mathbf{P}(X_0 = a_0 | X_j = b_j, j \leq -1)} - 1 \right| \leq \gamma_m. \quad (3)$$

We shall say that the process has *summable decay* if $\sum \gamma_m < +\infty$.

The next definition follows Ornstein (see, e.g., Ornstein, 1974).

Definition 2. The *canonical Markov approximation of order $k \in \mathbf{N}$* of a process $(X_n)_{n \in \mathbf{Z}}$ satisfying (1) is the stationary Markov chain of order k having as transition probabilities,

$$P^{(k)}(b | a_1, \dots, a_k) := \mathbf{P}(X_{k+1} = b | X_j = a_j, 1 \leq j \leq k) \quad (4)$$

for all integer $k \geq 1$ and $a_1, \dots, a_k, b \in A$.

We recall that a *coupling* of two processes $X = (X_n)_{n \in \mathbf{Z}}$ and $Y = (Y_n)_{n \in \mathbf{Z}}$ taking values in the alphabet A is any process $(\tilde{X}, \tilde{Y}) = (\tilde{X}_n, \tilde{Y}_n)_{n \in \mathbf{Z}}$ taking values in $A \times A$ such that $\text{Law}(\tilde{X}) = \text{Law}(X)$ and $\text{Law}(\tilde{Y}) = \text{Law}(Y)$.

Definition 3. The *distance \bar{d}* between two stationary processes X and Y is defined as

$$\bar{d}(X, Y) = \inf \left\{ \mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0) : (\tilde{X}, \tilde{Y}) \text{ stationary coupling of } X \text{ and } Y \right\}.$$

We now state our main result.

Theorem 4. Let $X = (X_n)_{n \in \mathbf{Z}}$ be a chain with complete connections and summable decay $(\gamma_m)_{m \geq 1}$. Then there is a constant $K > 0$ such that, for all $k \geq 1$,

$$\bar{d}(X, Y^{(k)}) \leq K \gamma_k,$$

where $Y^{(k)} = (Y_n^{(k)})_{n \in \mathbf{Z}}$ is the canonical Markov approximation of order k of the process X .

3. Proof of the theorem

The proof of the theorem is decomposed in the following way.

- First we introduce some notation.
- In Section 3.1, we prove a lemma showing that the transition probabilities of the approximating Markov chain are “close” to the transition probabilities of the original chain.
- In Section 3.2, we construct the coupling. We first define an appropriate system of transition probabilities \tilde{P} using the classical notion of maximal coupling (see Appendix A.1 in Barbour et al. (1992)). We then prove the existence of a stationary process $(\tilde{X}_n, \tilde{Y}_n)_{n \in \mathbb{Z}}$ with these transition probabilities.
- In Section 3.3, we obtain lower bounds for the probability of \tilde{X} being equal to \tilde{Y} during a certain number of steps given the history of the coupling. The more they have been equal in the past, the greater is this probability. If they were not equal at the previous step, they keep a positive (bounded away from 0) probability to become equal.
- The final estimation of $\mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0)$ is given in Section 3.4.

A sequence $\underline{x} = (x_j)_{j \leq -1}$ of elements of the alphabet A will be called a *history*. We shall denote by \underline{A} the set of all the histories. Given two histories \underline{x} and \underline{y} , the notation $\underline{x} \stackrel{m}{=} \underline{y}$ indicates that $x_j = y_j$ for all $-m \leq j \leq -1$. For the sake of notational simplicity, we shall denote

$$P(a | \underline{x}) = \mathbf{P}(X_0 = a | X_j = x_j, j \leq -1). \quad (5)$$

These objects exist for all $\underline{x} \in \underline{A}$ and $a \in A$ by virtue of Eq. (2). They admit three different interpretations. Firstly they can be seen as (a continuous version of) the conditional probabilities “knowing all the past” of the event $\{X_0 = a\}$. This motivates our notation. Secondly, they can be interpreted as transition probabilities that to each history associate (continuously) a law for the next step. Finally, one can think of them simply as functions from $A \times \underline{A}$ onto $[0, 1]$. Property (2) says that these functions are continuous while property (3) implies that they are indeed log-continuous. With notation (5), property (3) becomes

$$\sup \left\{ \left| \frac{P(a | \underline{x})}{P(a | \underline{y})} - 1 \right| ; \underline{x}, \underline{y} : \underline{x} \stackrel{m}{=} \underline{y} \right\} \leq \gamma_m, \quad (6)$$

with $a \in A$, $\underline{x}, \underline{y} \in \underline{A}$.

Let $P^{(k)}$ be the transition probability defined by (4). It is natural to use the same notation for the map from $A \times \underline{A}$ to $[0, 1]$ defined as

$$P^{(k)}(a | \underline{y}) = P^{(k)}(a | y_{-k}, \dots, y_{-1}). \quad (7)$$

With this notation, as soon as $\underline{x} \stackrel{k}{=} \underline{y}$, we have $P^{(k)}(a | \underline{y}) = P^{(k)}(a | \underline{x})$.

3.1. Properties of the Markov approximation

We now state the crucial consequences of property (6) for the transition probabilities of the canonical Markov approximation.

Lemma 5. For all integer $m > 0$,

$$\inf_{a \in A, \underline{x} \in \underline{A}} P^{(k)}(a | \underline{x}) \geq \inf_{a \in A, \underline{x} \in \underline{A}} P(a | \underline{x}) > 0, \quad (8)$$

$$\sup_{\underline{x}, \underline{y}: \underline{x} \stackrel{m}{=} \underline{y}} \sum_{a \in A} |P(a | \underline{x}) - P^{(k)}(a | \underline{y})| \leq \gamma_{m \wedge k}. \quad (9)$$

Proof. Property (6) guarantees that the functions $\underline{x} \rightarrow \log(P(a | \underline{x}))$ are continuous on the compact set \underline{A} . Hence, they are bounded for all a and the rightmost inequality in (8) follows. The conditional probability $\mathbf{P}(X_0 = a | X_j = y_j, -k \leq j \leq -1)$ can be written as an integral of $\underline{u} \rightarrow P(a | \underline{u})$ over a set on which $\underline{u} \stackrel{k}{=} \underline{y}$. Hence,

$$\inf_{\underline{u}: \underline{u} \stackrel{k}{=} \underline{y}} P(a | \underline{u}) \leq \mathbf{P}(X_0 = a | X_j = y_j, -k \leq j \leq -1) \leq \sup_{\underline{u}: \underline{u} \stackrel{k}{=} \underline{y}} P(a | \underline{u}). \quad (10)$$

It follows from Eqs. (4), (7) and (10), that

$$\inf_{\underline{u}: \underline{u} \stackrel{k}{=} \underline{y}} P(a | \underline{u}) \leq P^{(k)}(a | \underline{y}) \leq \sup_{\underline{u}: \underline{u} \stackrel{k}{=} \underline{y}} P(a | \underline{u}). \quad (11)$$

As $\inf_{\underline{u}: \underline{u} \stackrel{k}{=} \underline{x}} (P(a | \underline{u})) \geq \inf_{\underline{u} \in \underline{A}} (P(a | \underline{u}))$, this proves the leftmost inequality in (8).

Let us fix $a \in A$ and histories $\underline{x}, \underline{y}$ such that $\underline{x} \stackrel{m}{=} \underline{y}$ for some integer $m > 0$. According to Eq. (11), we have,

$$\left| \frac{P(a | \underline{x})}{P^{(k)}(a | \underline{y})} - 1 \right| \leq \sup \left\{ \left| \frac{P(a | \underline{u})}{P(a | \underline{v})} - 1 \right| ; \underline{u}, \underline{v} : \underline{u} \stackrel{k}{=} \underline{x}, \underline{v} \stackrel{k}{=} \underline{y} \right\}.$$

Noticing that $\underline{u} \stackrel{k}{=} \underline{x}, \underline{v} \stackrel{k}{=} \underline{y}$ and $\underline{x} \stackrel{m}{=} \underline{y}$ imply $\underline{u} \stackrel{k \wedge m}{=} \underline{v}$, and applying Eq. (6), we see that

$$|P(a | \underline{x}) - P^{(k)}(a | \underline{y})| \leq \gamma_{m \wedge k} P^{(k)}(a | \underline{y}).$$

We get Eq. (9) by summing over all the possible a . \square

Remark 6. In fact, Eq. (11) is the only property of the Markov transitions used in the sequel. Thus, our results apply to any Markov approximation scheme, not necessarily the canonical one, satisfying Eq. (11).

3.2. Construction of the coupling

We first define coupled transition probabilities. These are laws on $A \times A$ depending measurably on double histories, whose projections on each coordinate coincide, respectively, with the transition probabilities of the original and the approximating process. These transition probabilities are shown to be continuous and, hence, there exists a process compatible with them. This process is indeed a coupling of the original process and its canonical Markov approximation.

Given two 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 μ and ν (for more details see Appendix A.1 in Barbour et al. (1992)) that is a coupling which maximizes the weight of the diagonal. It can be defined as follows:

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

The important point here is that the distribution $\mu \tilde{\times} \nu$ on $A \times A$ satisfy simultaneously,

$$\begin{aligned}\sum_{a \in A} \mu \tilde{\times} \nu(a, a) &= \sum_{a \in A} \mu(a) \wedge \nu(a) \\ &= 1 - \sum_{a \in A} (\mu(a) - \nu(a))^+ \\ &= 1 - \frac{1}{2} \sum_{a \in A} |\mu(a) - \nu(a)|,\end{aligned} \tag{12}$$

$$\sum_{a \in A} \mu \tilde{\times} \nu(a, b) = \nu(b), \quad \sum_{b \in A} \mu \tilde{\times} \nu(a, b) = \mu(a).$$

Given the past, that is a double history $(\underline{x}, \underline{y})$, we set,

$$\tilde{P}((a, b) | (\underline{x}, \underline{y})) = P(\cdot | \underline{x}) \tilde{\times} P^{(k)}(\cdot | \underline{y})(a, b).$$

We now can state,

Proposition 7. *There is a stationary process $(\tilde{X}_n, \tilde{Y}_n)_{n \in \mathbb{Z}}$ taking values on $A \times A$ whose conditional probabilities satisfy,*

$$\mathbf{P}((\tilde{X}_0, \tilde{Y}_0) = (a, b) | (\tilde{X}_j, \tilde{Y}_j) = (x_j, y_j), j \leq -1) = \tilde{P}((a, b) | (\underline{x}, \underline{y})). \tag{13}$$

Moreover, under this probability, $\text{Law}(\tilde{X}) = \text{Law}(X)$ and $\text{Law}(\tilde{Y}) = \text{Law}(Y)$

Proof. We consider the functions \tilde{P} as a system of transition probabilities and we ask whether there exists a stationary process compatible with them. This is a rather classical problem. We notice that \tilde{P} depends continuously on P and $P^{(k)}$ which in turn depend continuously on $(a, b, \underline{x}, \underline{y})$. Hence, the transition probabilities \tilde{P} are continuous. A result by Ledrappier (1974) or Keane (1971) (concerning the so-called g -measures) proves the existence of a process satisfying Eq. (13).

Let $(\tilde{X}, \tilde{Y}) = (\tilde{X}_n, \tilde{Y}_n)_{n \in \mathbb{Z}}$ be such a process. Indeed, it appears from the construction that its marginal transition probabilities are what we need.

$$\mathbf{P}(\tilde{X}_0 = a | (\tilde{X}_j, \tilde{Y}_j) = (x_j, y_j), j \leq -1) = \sum_{b \in A} \tilde{P}((a, b) | (\underline{x}, \underline{y})) = P(a | \underline{x})$$

does not depend on \underline{y} . Hence, $\mathbf{P}(\tilde{X}_0 = a | \tilde{X}_j = x_j, j \leq -1) = P(a | \underline{x})$. The transition probabilities for \tilde{X} satisfy property (6) with summable decay. Ledrappier (1974) implies the unicity of the law of the processes compatible with these probabilities. As a consequence, $\text{Law}(\tilde{X}) = \text{Law}(X)$.

The proof that $\text{Law}(\tilde{Y}) = \text{Law}(Y)$ is even simpler: an analogous computation shows that $\mathbf{P}(\tilde{Y}_0 = b | \tilde{Y}_j = y_j, j \leq -1) = P^{(k)}(b | \underline{y})$. Hence \tilde{Y} is the only Markov chain compatible with the transition probabilities $P^{(k)}(b | y_{-k}, \dots, y_{-1})$.

Remark 8. The transition probabilities \tilde{P} do not define chains with complete connections because some of the transitions are zero. Moreover, in some situations, one can find arbitrarily close pairs of histories $(\underline{x}, \underline{y})$, $(\underline{x}', \underline{y}')$ such that $\tilde{P}((a, b) | (\underline{x}, \underline{y})) > 0$ but $\tilde{P}((a, b) | (\underline{x}', \underline{y}')) = 0$. Anyway, a direct computation proves that

$$|\tilde{P}((a, b) | (\underline{x}, \underline{y})) - \tilde{P}((a, b) | (\underline{x}', \underline{y}'))| \leq 4\gamma_m$$

holds, for all $a, b \in A$ and $\underline{x}, \underline{x}', \underline{y}, \underline{y}' \in \underline{A}$, with $\underline{x} \stackrel{m}{=} \underline{x}'$ and $\underline{y} \stackrel{m}{=} \underline{y}'$.

Let H be an event measurable with respect to the σ -algebra generated by $(\tilde{X}_n, \tilde{Y}_n)_{n \geq 0}$ and $(\underline{x}, \underline{y})$ a double history. From now on, we shall use the following short hand notation:

$$\mathbf{P}(H | (\underline{x}, \underline{y})) = \mathbf{P}(H | (\tilde{X}_j, \tilde{Y}_j) = (x_j, y_j), j \leq -1).$$

3.3. Main estimates

Let $\underline{x}, \underline{y}$ be two histories with $\underline{x} \stackrel{m}{=} \underline{y}$. We want to obtain an estimation of the probability of \tilde{X}_0 being different from \tilde{Y}_0 given these histories. First notice that, according to the consequence (12) of the definition of the coupling,

$$\mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0 | (\underline{x}, \underline{y})) = 1 - \sum_{a \in A} \tilde{P}((a, a) | (\underline{x}, \underline{y})) = \frac{1}{2} \sum_{a \in A} |P(a | \underline{x}) - P^{(k)}(a | \underline{y})|.$$

Let us define the sequence $(\tilde{\gamma}_n)_{n \in \mathbf{N}}$ by,

$$\tilde{\gamma}_0 = 1 - \inf_{a \in A, \underline{u} \in \underline{A}} P(a | \underline{u}),$$

$$\tilde{\gamma}_n = \min\left(\tilde{\gamma}_0, \frac{\gamma_n}{2}\right),$$

and let m_0 denote the first integer for which $\gamma_n \leq 2\tilde{\gamma}_0$. If $m \leq m_0$, we use Eq. (8), to see that,

$$\begin{aligned} \sum_a \tilde{P}((a, a) | (\underline{x}, \underline{y})) &\geq \inf_{a \in A} \tilde{P}((a, a) | (\underline{x}, \underline{y})) \\ &\geq \inf_{a \in A} (P(a | \underline{x}) \wedge P^{(k)}(a | \underline{y})) \\ &\geq \inf_{a \in A, \underline{u} \in \underline{A}} P(a | \underline{u}) \geq 1 - \tilde{\gamma}_0. \end{aligned}$$

If $m > m_0$ (provided $k > m_0$), we have, by Eq. (9),

$$\sum_{a \in A} |P(a | \underline{x}) - P^{(k)}(a | \underline{y})| \leq \gamma_{k \wedge m} \leq 2\tilde{\gamma}_{k \wedge m}.$$

We have that, for all $m \in \mathbf{N}$, and for all histories $\underline{x}, \underline{y}$ with $\underline{x} \stackrel{m}{=} \underline{y}$,

$$\mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0 | (\underline{x}, \underline{y})) \leq \tilde{\gamma}_{k \wedge m}. \quad (14)$$

Let us denote by $\Delta_{m,n}$ the sets $\Delta_{m,n} := \bigcap_{p=m}^n \{\tilde{X}_j = \tilde{Y}_j\} = \{\tilde{X}_j = \tilde{Y}_j, m \leq j \leq n\}$ and by $\Delta_{m,n}^c$ their complementary sets. Notice that $\Delta_{-m,-1}$ is the reunion over all the sequences $\underline{x}, \underline{y}$ with $\underline{x} \stackrel{m}{=} \underline{y}$ of the events $\{(\tilde{X}_j, \tilde{Y}_j) = (x_j, y_j); j \leq -1\}$.

Lemma 9. For all integers m, n and all double histories $(\underline{x}, \underline{y})$ with $\underline{x} \stackrel{m}{=} \underline{y}$,

$$\mathbf{P}(\Delta_{0,n} | (\underline{x}, \underline{y})) \geq \prod_{p=0}^n (1 - \tilde{\gamma}_{k \wedge (m+p)}). \quad (15)$$

Proof. Let $\underline{x}, \underline{y}$ be two histories with $\underline{x} \stackrel{m}{=} \underline{y}$. We write,

$$\begin{aligned} \mathbf{P}(\Delta_{0,n} | (\underline{x}, \underline{y})) &= \mathbf{P}(\tilde{X}_0 = \tilde{Y}_0 | (\underline{x}, \underline{y})) \prod_{p=1}^n \mathbf{P}(\tilde{X}_p = \tilde{Y}_p | \Delta_{0,p-1}, (\underline{x}, \underline{y})) \\ &= (1 - \mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0 | (\underline{x}, \underline{y}))) \prod_{p=1}^n (1 - \mathbf{P}(\tilde{X}_p \neq \tilde{Y}_p | \Delta_{0,p-1}, (\underline{x}, \underline{y}))) \\ &= \prod_{p=0}^n (1 - \mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0 | H_{m+p}^{(\underline{x}, \underline{y})})), \end{aligned} \quad (16)$$

where $H_{m+p}^{(\underline{x}, \underline{y})}$ is the event corresponding to the set of double histories $(\underline{u}, \underline{v})$ with $\underline{u} \stackrel{p}{=} \underline{v}$ and $u_{-p+j} = x_j, v_{-p+j} = y_j$ for all $j \leq -1$. Notice that $\underline{u} \stackrel{m+p}{=} \underline{v}$ for all histories $(\underline{u}, \underline{v})$ corresponding to an element of $H_{m+p}^{(\underline{x}, \underline{y})}$. That is, $H_{m+p}^{(\underline{x}, \underline{y})} \subset \Delta_{-m-p, -1}$. Using the same kind of arguments that yield inequality (10), we see that,

$$\mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0 | H_{m+p}^{(\underline{x}, \underline{y})}) \leq \sup_{(\underline{u}, \underline{v}) \in H_{m+p}^{(\underline{x}, \underline{y})}} \mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0 | (\underline{u}, \underline{v})) \leq \sup_{\underline{u} \stackrel{m+p}{=} \underline{v}} \mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0 | (\underline{u}, \underline{v})).$$

The lemma follows from this, (14) and (16).

From this result, we easily deduce,

Lemma 10.

$$\mathbf{P}(\Delta_{0,k-1} | \Delta_{-k,-1}) \geq (1 - \tilde{\gamma}_k)^k, \quad (17)$$

$$\mathbf{P}(\Delta_{0,k-1} | \Delta_{-k,-1}^c) \geq \prod_{p=0}^{+\infty} (1 - \tilde{\gamma}_p). \quad (18)$$

Proof. Using again the arguments yielding inequality (10), we have, for $H = \Delta_{-k,-1}$ and for $H = \Delta_{-k,-1}^c$,

$$\mathbf{P}(\Delta_{0,k-1} | H) \geq \inf_{(\underline{x}, \underline{y}) \in H} \mathbf{P}(\Delta_{0,k-1} | (\underline{x}, \underline{y})).$$

Hence, using Lemma 9 for $n = k - 1, m = k$, we obtain,

$$\mathbf{P}(\Delta_{0,k-1} | \Delta_{-k,-1}) \geq \prod_{p=0}^{k-1} (1 - \tilde{\gamma}_{k \wedge (k+p)}) = \prod_{p=0}^{k-1} (1 - \tilde{\gamma}_k) = (1 - \tilde{\gamma}_k)^k,$$

and, using Lemma 9 for $n = k - 1$, $m = 0$,

$$\mathbf{P}(\Delta_{0,k-1} | \Delta_{-k,-1}^c) \geq \prod_{p=0}^{k-1} (1 - \tilde{\gamma}_{k \wedge p}) \geq \prod_{p=0}^{k-1} (1 - \tilde{\gamma}_p) \geq \prod_{p=0}^{+\infty} (1 - \tilde{\gamma}_p).$$

Lemma 11.

$$\mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0) \leq \frac{1}{\prod_{m=0}^{+\infty} (1 - \tilde{\gamma}_m)} \frac{\mathbf{P}(\Delta_{0,k-1}^c)}{k}. \quad (19)$$

Proof. For all finite family $(A_i)_{i=1..k}$ of measurable sets, we have the decomposition

$$\bigcup_{i=1}^k A_i = \bigcup_{i=1}^k \left(A_i \setminus \left(A_i \cap \left(\bigcup_{j=i+1}^k A_j \right) \right) \right).$$

Notice that the last element of this partition is exactly A_k . Hence,

$$\begin{aligned} \mathbf{P}\left(\bigcup_{i=1}^k A_i\right) &= \mathbf{P}\left(\bigcup_{i=1}^k \left(A_i \setminus \left(A_i \cap \left(\bigcup_{j=i+1}^k A_j \right) \right) \right)\right) \\ &= \sum_{i=1}^k \mathbf{P}\left(A_i \setminus \left(A_i \cap \left(\bigcup_{j=i+1}^k A_j \right) \right) \right) \\ &= \sum_{i=1}^k \mathbf{P}(A_i) - \sum_{i=1}^{k-1} \mathbf{P}\left(A_i \cap \left(\bigcup_{j=i+1}^k A_j \right) \right). \end{aligned}$$

We use this decomposition to compute the probability of $\Delta_{i,k-1}^c = \bigcup_{j=i}^{k-1} \{\tilde{X}_j \neq \tilde{Y}_j\}$,

$$\begin{aligned} \mathbf{P}(\Delta_{0,k-1}^c) &= \sum_{i=0}^{k-1} \mathbf{P}(\tilde{X}_i \neq \tilde{Y}_i) - \sum_{i=0}^{k-2} \mathbf{P}(\{\tilde{X}_i \neq \tilde{Y}_i\} \cap \Delta_{i+1,k-1}^c) \\ &= \sum_{i=0}^{k-1} \mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0) - \sum_{i=0}^{k-2} \mathbf{P}(\Delta_{i+1,k-1}^c | \tilde{X}_i \neq \tilde{Y}_i) \mathbf{P}(\tilde{X}_i \neq \tilde{Y}_i) \\ &= k \mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0) - \sum_{i=0}^{k-2} \mathbf{P}(\Delta_{0,k-i-2}^c | \tilde{X}_{-1} \neq \tilde{Y}_{-1}) \mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0). \end{aligned} \quad (20)$$

Let us now notice that, according to Lemma 9,

$$\mathbf{P}(\Delta_{0,k-i-2} | \tilde{X}_{-1} \neq \tilde{Y}_{-1}) \geq \prod_{m=0}^{k-i-1} (1 - \tilde{\gamma}_{k \wedge m}) \geq \prod_{m=0}^{+\infty} (1 - \tilde{\gamma}_m). \quad (21)$$

Inequalities (20) and (21) yield the lemma.

3.4. Conclusion of the proof

We now have all the elements to prove Theorem 4. From

$$\begin{aligned} \mathbf{P}(\Delta_{0,k-1}^c) &= \mathbf{P}(\Delta_{0,k-1}^c | \Delta_{-k,-1}) \mathbf{P}(\Delta_{-k,-1}) + \mathbf{P}(\Delta_{0,k-1}^c | \Delta_{-k,-1}^c) \mathbf{P}(\Delta_{-k,-1}^c) \\ &\leq (1 - (1 - \tilde{\gamma}_k)^k) + \left(1 - \prod_{p=0}^{+\infty} (1 - \tilde{\gamma}_p)\right) \mathbf{P}(\Delta_{0,k-1}^c), \end{aligned}$$

we deduce that

$$\mathbf{P}(\Delta_{0,k-1}^c) \leq \frac{1 - (1 - \tilde{\gamma}_k)^k}{\prod_{p=0}^{+\infty} (1 - \tilde{\gamma}_p)}. \quad (22)$$

Using Lemma 11, we get,

$$\begin{aligned} \mathbf{P}(\tilde{X}_0 \neq \tilde{Y}_0) &\leq \frac{1}{k \prod_{p=0}^{+\infty} (1 - \tilde{\gamma}_p)} \mathbf{P}(\Delta_{0,k-1}^c) \\ &\leq \frac{1}{(\prod_{p=0}^{+\infty} (1 - \tilde{\gamma}_p))^2} \frac{1 - (1 - \tilde{\gamma}_k)^k}{k}. \end{aligned}$$

To conclude the proof we notice that, on the one hand,

$$1 - (1 - \tilde{\gamma}_k)^k \sim 1 - e^{k \log(1 - \frac{\gamma_k}{2})} \sim 1 - e^{-\frac{k}{2} \gamma_k} \sim \frac{k}{2} \gamma_k,$$

because, as $(\gamma_m)_{m \geq 0}$ is decreasing and summable, $k\gamma_k \rightarrow 0$, and, on the other hand,

$$\prod_{p=0}^{+\infty} (1 - \tilde{\gamma}_p) > 0,$$

because, $\log \prod_{p=0}^n (1 - \tilde{\gamma}_p) = \sum_{p=0}^n \log(1 - \tilde{\gamma}_p) \sim -\frac{1}{2} \sum_{p=0}^n \gamma_p$ and $\sum_{p=0}^{+\infty} \gamma_p < +\infty$.

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