

# On Categorical Time Series Models With Covariates

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## Abstract

We study the problem of stationarity and ergodicity for autoregressive multinomial logistic time series models which possibly include a latent process and are defined by a GARCH-type recursive equation. We improve considerably upon the existing results related to stationarity and ergodicity conditions of such models. Proofs are based on theory developed for chains with complete connections. This approach is based on a useful coupling technique which is utilized for studying ergodicity of more general finite-state stochastic processes. Such processes generalize finite-state Markov chains by assuming infinite order models of past values. For finite order Markov chains, we also discuss ergodicity properties when some strongly exogenous covariates are considered in the dynamics of the process.

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# 1 Introduction

The goal of this article is to improve upon theoretical properties of regression based models for the analysis of categorical time series that might include some covariates. Binary time series are particular cases of a categorical time series and the results we obtain apply to this case as well. We take the point of view of generalized linear models theory; see [McCullagh and Nelder \(1989\)](#). The conditional distribution of a categorical time series given its past is multinomial which obviously belongs to the multivariate exponential family of distributions. As such, the theory of generalized linear models can be applied for modeling different types of categorical data; nominal, interval and scale. We will be mostly concerned with nominal data and therefore the multinomial logistic model is the natural candidate for model fitting; see [Fahrmeir and Tutz \(2001\)](#) and [Kedem and Fokianos \(2002\)](#), among other references, for further discussion on modeling issues regarding categorical data. We emphasize that finite state Markov chains provide a simple but prominent model of a categorical time series where lagged values of the response affect the determination of its future states. However, Markov modeling in the context of categorical time series, poses challenging problems. Indeed, as the order of the Markov chain increases so does the number of free parameters; in fact, the number of free parameters increases exponentially fast. Furthermore, the Markovian property requires simultaneous specification of the dynamics of the response and any possible covariates observed jointly; such a specification might not be possible, in general.

We will be studying models for binary and, more generally, categorical time series, which are of infinite order or they are driven by a latent process or a feedback mechanism. This type of models is quite analogous to GARCH models -see [Bollerslev \(1986\)](#)- but they are defined in terms of conditional log-odds instead of conditional variances. In particular, feedback models make possible low dimensional parametrization, yet they can accommodate quite complicated data structures. Examples of feedback models, in the context of binary and categorical time series have been studied recently by [Moysiadis and Fokianos \(2014\)](#) and [Fokianos and Moysiadis \(2017\)](#), among others. We will discuss these results and we will compare them with our findings which improve these works. Models and inference about binary time series, in general, are topics that have been studied by several authors; see [Kedem \(1980\)](#) for an early treatment. Regression modeling, in this context, has been studied by [Cox \(1981\)](#), [Stern and Coe \(1984\)](#), and [Slud and Kedem \(1994\)](#), among others; see also [Kedem and Fokianos \(2002, Ch. 2-3\)](#) for other early references. Recently, binary time series data have been increasingly popular in various financial applications ([Breen et al. \(1989\)](#), [Butler and Malaikah \(1992\)](#), [Christoffersen and Diebold \(2006\)](#), [Christoffersen et al. \(2007\)](#), [Startz \(2008\)](#),

Nyberg (2010, 2011, 2013), Kauppi (2012) and Wu and Cui (2014)), but also to other scientific fields. Previous results related to theoretical properties of such models were given by de Jong and Woutersen (2011).

Related work on categorical time series has been reported by Fahrmeir and Kaufmann (1987), Kaufmann (1987), Fokianos and Kedem (2003) and Russell and Engle (1998, 2005) who proposed a categorical time series model for financial transactions data. Alternative classes of models are based on the probit link function. Such autoregressive models have been considered by Zeger and Qaqish (1988), Rydberg and Shephard (2003), Kauppi and Saikkonen (2008), among others. Several other classes of models for the analysis of categorical data have been studied; see the books by Joe (1997) and MacDonald and Zucchini (1997) and the articles by Biswas and Song (2009) and Weiß (2011).

To prove the theoretical results, we will be assuming a contraction type condition; such conditions are usually employed for the theoretical analysis of time series models. For instance, in the case of count time series models, see Doukhan et al. (2012), Fokianos et al. (2009) and Neumann (2011). However, our work is closely related to the modeling approach suggested by Fokianos and Tjøstheim (2011), because the main idea is essentially to employ the so called canonical link process to model the observed data. Note that de Jong and Woutersen (2011) have shown near epoch dependence for a binary time series models but these authors have a different modeling point of view.

Likelihood based inference for the models we study can be developed along the lines of previous references. The proof of consistency and asymptotic normality is based on standard arguments concerning convergence of the score function and the Hessian matrix. However, we mention that this work relaxes considerably previous results. for the case of a model with covariates we improve upon Kaufmann (1987) and Kedem and Fokianos (2002, Ch.3) because we avoid any assumptions on the design of covariates or their boundedness. Since proofs of these facts have been documented in several of the previous references, we do not give any details.

The article is structured as follows: Section 2 discusses general categorical time series models by allowing the conditional probabilities to depend on the whole past of the series. In addition we will be giving a result about the stationarity and ergodicity of chains with complete connections. These results this will be applied to the case of an infinite order autoregressive multinomial logistic model. Section 3 discuss models which might include a latent process. The results obtained by Theorem 1 are applicable to the case of models considered by Moysiadis and Fokianos (2014) and Fokianos and Moysiadis (2017) and improve the

previously given stationarity and ergodicity conditions. Finally, Section 4 discuss inclusion of exogenous covariates to the autoregressive multinomial logistic model; Theorem 2 is the main result of this section which discuss existence of such processes and their ergodic properties.

## 2 Time series autoregressive models for categorical data

### 2.1 A general approach

Let  $A$  be a finite set. For simplicity, we assume that  $A = \{1, 2, \dots, N\}$ , where  $N$  is a nonnegative integer. Suppose that we observe a process with state space  $A$  and we are interested on modeling its dynamics. For instance, consider modeling of a stock price change (0 for no change, 1 for positive change and -1 for a negative change; see [Russell and Engle \(1998\)](#)) or sleep state status (see [Fokianos and Kedem \(2003\)](#)). Towards this goal, define a  $(N - 1)$ -dimensional vector  $Y_t = (Y_{1t}, Y_{2t}, \dots, Y_{(N-1)t})'$ , for  $1 \leq t \leq n$ , such that

$$Y_{jt} = \begin{cases} 1, & \text{if the } j\text{'th category is observed at time } t, \\ 0, & \text{otherwise,} \end{cases}$$

for all  $j = 1, 2, \dots, N-1$ . Throughout this work, consider a stochastic processes  $(Y_t)_{t \in \mathbb{Z}}$  adapted to a filtration  $(\mathcal{F}_t)_{t \in \mathbb{Z}}$  which is defined through a vector of conditional "success" probabilities, say  $p_t \equiv (p_{1t}, p_{2t}, \dots, p_{(N-1)t})'$ .

In other words

$$p_{jt} = \mathbb{P}(Y_{jt} = 1 | \mathcal{F}_{t-1}), \quad 1 \leq j \leq N-1. \quad (1)$$

For the last category  $N$ , set  $Y_{Nt} = 1 - \sum_{j=1}^{N-1} Y_{jt}$  and corresponding success probability  $p_{Nt} = 1 - \sum_{j=1}^{N-1} p_{jt}$ .

There are several possibilities for autoregressive modeling of processes that take values on a finite space. For instance, assuming that  $d$  is a vector and  $A, B$  matrices of appropriate dimension, consider the following linear model

$$p_t = d + Ap_{t-1} + BY_{t-1}, \quad t \in \mathbb{Z}, \quad (2)$$

which was studied by [Russell and Engle \(1998\)](#) and [Qaqish \(2003\)](#). Model (2) implies quite complex restrictions on the parameters  $d, A$  and  $B$  because each element of the vector  $p_t$  has to belong in the interval  $(0, 1)$ . Such restrictions become even more involved when a covariate process is included in (2). To avoid such subtle technicalities, we adapt the generalized linear models point of view; see [Fokianos and Kedem \(2003\)](#) for instance. For  $j = 1, 2, \dots, N-1$ , define

$$\lambda_{jt} = \log(p_{jt}/p_{Nt})$$

and suppose that the vector process  $\lambda_t = (\lambda_{1t}, \dots, \lambda_{(N-1)t})'$  is determined by the infinite order model

$$\lambda_t = g(Y_{t-1}, Y_{t-2}, \dots), \quad (3)$$

where  $g$  is a suitably defined function. Then, the process  $(Y_t)_{t \in \mathbb{Z}}$  which satisfies (1) and (3), takes its values in the set  $E = \{e_1, e_2, \dots, e_{N-1}, \mathbf{0}\}$  where  $\{e_1, \dots, e_{N-1}\}$  is the canonical basis of  $\mathbb{R}^{N-1}$  and  $\mathbf{0}$  is the null vector of  $\mathbb{R}^{N-1}$ . Furthermore,  $g : E^{\mathbb{N}} \rightarrow \mathbb{R}^{N-1}$  is a measurable function and the conditional distribution of  $Y_t$  given its past values  $Y_{t-1}^- \equiv (Y_{t-1}, Y_{t-2}, \dots)$  possibly depends on its infinite past. A useful example of such process is given by the linear process

$$\lambda_t = d + \sum_{j \geq 1} A_j Y_{t-j} \quad (4)$$

where  $d$  is a  $(N-1)$ -dimensional vector and  $(A_j)_{j \geq 1}$  is a sequence of  $(N-1) \times (N-1)$  matrices. Comparison of (4) to (2) shows that unnecessary restrictions on the unknown coefficients can be circumvented since the vector  $\lambda_t \in \mathbb{R}^{N-1}$ . Furthermore, covariates can be easily included in (4) by including an additional additive term. Other categorical type autoregressive models can be considered but (4) has been widely used in several applications. In the case that  $N = 2$ , then (4) is a simple logistic regression model which has been studied widely in the literature (see [Cox and Snell \(1970\)](#) for an early reference).

Processes, as those we consider in this work, are particular examples of a more general class of processes which are called *chains with complete connections*. Such processes have been widely studied in applied probability; [Doebelin and Fortet \(1937\)](#), [Harris \(1955\)](#) and [Iosifescu and Grigorescu \(1990\)](#). Following the work of [Bressaud et al. \(1999\)](#), we discuss next a coupling technique related to chains with complete connections.

## 2.2 Some results about chains with complete connection

Throughout this section, consider a finite state space  $E$ . For  $x, y \in E^{\mathbb{N}}$  and a positive integer  $m$ , we write  $x \stackrel{m}{=} y$  if  $x_i = y_i$  for  $0 \leq i \leq m-1$ . Consider a probability kernel  $p(\cdot | \cdot)$  defined on  $(E^{\mathbb{N}}, \mathcal{B}(E^{\mathbb{N}}))$  and takes values on  $(E, \mathcal{B}(E))$  which satisfies the following assumption:

**Assumption (A)** There exists a sequence  $(\gamma_m)_{m \in \mathbb{N}}$  which decreases to zero, as  $m \rightarrow \infty$ , with  $\gamma_0 < 1$  and such that for  $a \in E$

$$\inf_{x, y: x \stackrel{m}{=} y} \frac{p(a|x)}{p(a|y)} \geq 1 - \gamma_m.$$

A chain with complete connections is a stationary process satisfying Assumption (A).

For  $x \in E^{\mathbb{N}}$ , consider the chain  $(Z_n^x)_{n \in \mathbb{Z}}$  which satisfies that  $Z_{-j}^x = x_j$  for  $j \geq 1$  and

$$\mathbb{P}(Z_n^x = a | Z_{n-j}^x = z_j, j \geq 1) = p(a|z) \prod_{j=n+1}^{\infty} \mathbb{1}_{z_j = x_{j-n}}.$$

In addition, given a real-valued sequence  $(\gamma_n)_{n \in \mathbb{N}}$ , let the Markov chain  $(S_n^{(\gamma)})_{n \in \mathbb{N}}$  taking values in  $\mathbb{N}$  and defined by

$$\mathbb{P}(S_0^{(\gamma)} = 0) = 1, \quad \mathbb{P}(S_{n+1}^{(\gamma)} = i + 1 | S_n^{(\gamma)} = i) = 1 - \gamma_i, \quad \mathbb{P}(S_{n+1}^{(\gamma)} = 0 | S_n^{(\gamma)} = i) = \gamma_i.$$

For  $n \geq 1$ , define the quantity

$$\gamma_n^* = \mathbb{P}(S_n^{(\gamma)} = 0),$$

which plays a crucial rule for evaluating the mixing coefficients of the chain. The following result is given by [Bressaud et al. \(1999, Prop. 1 and Lemma 1\)](#).

**Proposition 1.** *For all  $x, y \in E^{\mathbb{N}}$ , there is a coupling  $((U_n^{x,y}, V_n^{x,y}))_{n \in \mathbb{Z}}$  of  $(Z_n^x)_{n \in \mathbb{Z}}$  and  $(Z_n^y)_{n \in \mathbb{Z}}$  such that the integer-valued process  $(T_n^{x,y})_{n \in \mathbb{Z}}$  defined by*

$$T_n^{x,y} = \inf \{m \geq 0 : U_{n-m}^{x,y} \neq V_{n-m}^{x,y}\},$$

satisfies

$$\mathbb{P}(S_n^{(\gamma)} \geq k) \leq \mathbb{P}(T_n^{x,y} \geq k) \quad \forall k \in \mathbb{N}.$$

Proposition 1 is proved by defining iteratively the pair  $(U_n^{x,y}, V_n^{x,y})$  using the maximal coupling of the conditional distributions  $p(\cdot | (u_{n-j})_{j \geq 1})$  and  $p(\cdot | (v_{n-j})_{j \geq 1})$  (i.e the coupling associated to the total variation distance between these conditional distributions). Proposition 1 yields the following corollary (see also [Bressaud et al. \(1999, Cor. 1\)](#) for a specific case of the following result).

**Corollary 1.** *For all  $k \geq 1$ ,  $x, y \in E^{\mathbb{N}}$  and  $B \in \mathcal{B}(E^k)$ , we have*

$$\left| \mathbb{P}((Z_n^x, \dots, Z_{n+k}^x) \in B) - \mathbb{P}((Z_n^y, \dots, Z_{n+k}^y) \in B) \right| \leq \sum_{j=0}^k \left( \prod_{m=0}^{j-1} (1 - \gamma_m) \right) \gamma_{n+k-j}^*.$$

*Proof.* Using Proposition 1, we obtain

$$\begin{aligned} \left| \mathbb{P}((Z_n^x, \dots, Z_{n+k}^x) \in B) - \mathbb{P}((Z_n^y, \dots, Z_{n+k}^y) \in B) \right| &\leq \mathbb{P}((U_n^{x,y}, \dots, U_{n+k}^{x,y}) \neq (V_n^{x,y}, \dots, V_{n+k}^{x,y})) \\ &\leq \mathbb{P}(T_n^{x,y} \leq k). \end{aligned}$$

Proposition 1 implies that  $\mathbb{P}(T_n^{x,y} \leq k) \leq \mathbb{P}(S_n^{(\gamma)} \leq k)$ . The result of the corollary now follows by bounding  $\mathbb{P}(S_n^{(\gamma)} \leq k)$  along the lines of the derivation of Bressaud et al. (1999, eq. (4.25)).  $\square$

As pointed out in Bressaud et al. (1999), if  $\lim_{n \rightarrow \infty} \gamma_n^* = 0$ , Corollary 1 implies existence and uniqueness of a stationary chain  $(Z_n)_{n \in \mathbb{Z}}$  with complete connections and satisfying Assumption (A).

Furthermore, Corollary 1 yields a bound for controlling the  $\phi$ -mixing coefficients associated with  $(Z_n)_{n \in \mathbb{Z}}$ . Indeed, recall that for two  $\sigma$ -algebras  $\mathcal{A}$  and  $\mathcal{B}$ , their  $\phi$ -mixing coefficients are defined by (see Doukhan (1994), for instance)

$$\phi(\mathcal{A}, \mathcal{B}) = \sup_{(A, B) \in \mathcal{A} \times \mathcal{B}: \mathbb{P}(A) > 0} (|\mathbb{P}(B|A) - \mathbb{P}(B)|).$$

Then, the  $\phi$ -mixing coefficients of the random sequence  $(Z_n)_{n \in \mathbb{Z}}$  are given by

$$\phi(n) = \phi(\mathcal{F}_{-\infty, 0}, \mathcal{F}_{n, \infty}).$$

Since a Borel set on the infinite product can be approximated by a finite union of cylinder sets, we have also

$$\phi(n) = \sup_{k \in \mathbb{N}} \phi(\mathcal{F}_{-\infty, 0}, \mathcal{F}_{n, n+k}).$$

**Proposition 2.** *Suppose that  $\sum_{n \geq 1} \gamma_n^* < \infty$ . Then the infinite order stationary Markov chain  $(Z_n)_{n \in \mathbb{Z}}$ , which exists by Corollary 1, is  $\phi$ -mixing with mixing coefficients satisfying  $\phi(n) \leq \sum_{j \geq n} \gamma_j^*$ .*

*Proof.* Suppose that  $\mu$  denotes the probability distribution of  $(Z_i)_{i \leq -1}$ . Then

$$\begin{aligned} \left| \mathbb{P}((Z_n^x, \dots, Z_{n+k}^x) \in B) - \mathbb{P}((Z_n, \dots, Z_{n+k}) \in B) \right| &\leq \int \left| \mathbb{P}((Z_n^x, \dots, Z_{n+k}^x) \in B) - \mathbb{P}((Z_n^y, \dots, Z_{n+k}^y) \in B) \right| \mu(dy) \\ &\leq \sum_{j=0}^k \left( \prod_{m=0}^{j-1} (1 - \gamma_m) \right) \gamma_{n+k-j}^*. \end{aligned}$$

But the last bound does not depend on  $x$ . Hence, we obtain

$$\begin{aligned} \left| \mathbb{P}((Z_n, \dots, Z_{n+k}) \in B | \mathcal{F}_{-\infty, 0}) - \mathbb{P}((Z_n, \dots, Z_{n+k}) \in B) \right| &\leq \sum_{j=0}^k \left( \prod_{m=0}^{j-1} (1 - \gamma_m) \right) \gamma_{n+k-j}^* \\ &\leq \sum_{j \geq n} \gamma_j^*. \end{aligned}$$

The last bound, which does not depend on  $k$  and  $B$ , is also an upper bound for  $\phi(n)$ .  $\square$

**Remark 1.** It has been shown in Bressaud et al. (1999, Prop. 2) that

$$\sum_{k \geq 1} \gamma_k < \infty \Rightarrow \sum_{k \geq 1} \gamma_k^* < \infty.$$

Moreover, if  $(\gamma_m)_m$  decreases exponentially, then so does  $(\gamma_n^*)_n$ . Hence, the result of Prop. 2 follows again and if  $(\gamma_m)_m$  decreases to zero exponentially fast then so does  $(\phi(n))_n$ . Note also that the  $\phi$ -mixing property implies ergodicity of the process; see Bradley (2007, pp. 50-51).

### 2.3 Application to categorical time series

Recall the categorical time series model  $(Y_t)_{t \in \mathbb{Z}}$  whose state space is  $E = \{e_1, \dots, e_{N-1}, \mathbf{0}\}$  and defined by (1) and (3). From the results of the previous subsection, we deduce the following corollary. ( $\|\cdot\|$  denotes the Euclidian norm on  $\mathbb{R}^{N-1}$ .)

**Corollary 2.** Assume model (3) and let a function  $g : E^{\mathbb{N}} \rightarrow \mathbb{R}^{N-1}$  be such that there exist a sequence  $(\delta_j)_{j \in \mathbb{N}}$  which satisfies  $\sum_{j \in \mathbb{N}} \sum_{k \geq j} \delta_k < \infty$  and

$$\|g(x) - g(y)\| \leq \sum_{j \in \mathbb{N}} \delta_j \mathbb{1}_{x_j \neq y_j}. \quad (5)$$

Then, there exists a unique stochastic process  $(Y_t)_{t \in \mathbb{Z}}$  taking values in  $E$  such that

$$\mathbb{P}(Y_t = e_j | \mathcal{F}_{t-1}) = \frac{\exp(g_j(Y_{t-1}, Y_{t-2}, \dots))}{1 + \sum_{s=1}^{N-1} \exp(g_s(Y_{t-1}, Y_{t-2}, \dots))}, \quad 1 \leq j \leq N-1. \quad (6)$$

Moreover  $(Y_t)_{t \in \mathbb{Z}}$  is stationary and  $\phi$ -mixing.

*Proof.* Denote by  $p(\cdot|\cdot)$  the probability kernel defined by

$$p(e_j|x) = F_j[g(x_0, x_1, \dots)], \quad 1 \leq j \leq N-1,$$

where  $F_j : \mathbb{R}^{N-1} \rightarrow [0, 1]$  is defined for  $z \in \mathbb{R}^{N-1}$  by

$$F_j(z) = \frac{\exp(z_j)}{1 + \sum_{s=1}^{N-1} \exp(z_s)}, \quad 1 \leq j \leq N-1.$$

Because of (1),  $F_N(z) = \left(1 + \sum_{s=1}^{N-1} \exp(z_s)\right)^{-1}$ . The Lipschitz assumption (5) implies that the  $j$ 'th component of  $g$ , say  $g_j$ , is bounded, for  $j = 1, 2, \dots, N$ . Hence, there exists  $\eta > 0$  such that, for all  $1 \leq j \leq N-1$  and  $x \in E^{\mathbb{N}}$ ,

$$\eta \leq p(e_j|x), \quad \eta \leq p(\mathbf{0}|x).$$

Moreover,  $F'_j$  is bounded, for all  $j$ . Set  $M = \max_{1 \leq j \leq N-1} \sup_{z \in \mathbb{R}^{N-1}} \|F'_j(z)\|$ . Then, if  $x \stackrel{m}{=} y$  and  $a \in E$ , we have that

$$\frac{p(a|x)}{p(a|y)} \geq 1 - \frac{M \sum_{j \geq m+1} \delta_j}{\eta}.$$



Provided that  $m$  is large enough, choose  $\gamma_m = M \sum_{j \geq m+1} \delta_j / \eta$ . Hence, there exists an  $m$  such that  $\gamma_m = 1 - \eta$ . Then we have  $\sum_{k \geq 1} \gamma_k < \infty$  and using Remark 1, we have also  $\sum_{k \geq 1} \gamma_k^* < \infty$ . Then from Corollary 1 and Proposition 2, there exists a unique stationary solution  $(Y_t)_{t \in \mathbb{Z}}$  satisfying (6) and the solution is  $\phi$ -mixing.  $\square$

We note that the condition  $\sum_{j \in \mathbb{N}} \sum_{k \geq j} a_k < \infty$  is equivalent to the condition  $\sum_{j \in \mathbb{N}} j a_j < \infty$ . For the *infinite order* linear model (4), Corollary 2 applies provided that  $\sum_{j \geq 1} j \|A_j\| < \infty$  where  $\|A_j\|$  denotes the corresponding operator norm of the matrix  $A_j$ . In particular, when  $N = 2$ , we obtain that the logistic autoregressive model of infinite order is stationary and  $\phi$ -mixing if  $\sum_{j \geq 1} j |A_j| < \infty$ , where  $(A_j)_{j \geq 1}$  belongs to  $\mathbb{R}$ .

### 3 Categorical time series with a latent process

In this section, we consider some specific instances of chains with complete connections. Following the methodology of GARCH models (see Engle (1982), Bollerslev (1986) and the book by Francq and Zakoian (2010) for instance), and recalling the notation introduced in (3) we model the latent process  $(\lambda_t)_{t \in \mathbb{Z}}$  to depend additionally on its past values. From a statistical perspective, such parametrization yields parsimony and allows for more flexible structures that can accommodate various forms of autocorrelation. To be specific, suppose that  $p$  and  $q$  are two positive integers and Let  $f : \mathbb{R}^{(N-1)p} \times E^q \rightarrow \mathbb{R}^{N-1}$  be a function such that

$$\lambda_t = f(\lambda_{t-1}, \dots, \lambda_{t-p}, Y_{t-1}, \dots, Y_{t-q}), \quad t \in \mathbb{Z}. \quad (7)$$

We will say that the process  $((Y_t, \lambda_t))_{t \in \mathbb{Z}}$  is a solution of the problem  $\mathcal{P}_f$  if (7) is satisfied and for each  $t \in \mathbb{Z}$ ,  $\lambda_t$  is  $\mathcal{F}_{t-1}$ -measurable.

#### 3.1 A general result

For  $\underline{y} \in E^q$ , define the mapping  $G_{\underline{y}} : \mathbb{R}^{(N-1)p} \rightarrow \mathbb{R}^{(N-1)p}$  by

$$G_{\underline{y}}(\underline{x}) = \left( f(x_1, \dots, x_p, y_1, \dots, y_q)', x'_1, \dots, x'_{p-1} \right)',$$

where  $f(\cdot)$  has been defined by (7). The main result of this section is the following.

**Theorem 1.** Suppose that there exist an integer  $k \geq 1$ ,  $\kappa \in (0, 1)$  and  $K > 0$  such that

$$\|G_{\underline{y}}(\underline{x}) - G_{\underline{y}'}(\underline{x}')\| \leq K(\|\underline{x} - \underline{x}'\| + \mathbb{1}_{\underline{y} \neq \underline{y}'}),$$

and for all  $\underline{x}, \underline{x}', \underline{y}_1, \dots, \underline{y}_q$ ,

$$\|G_{\underline{y}_1} \circ \dots \circ G_{\underline{y}_k}(\underline{x}) - G_{\underline{y}_1} \circ \dots \circ G_{\underline{y}_k}(\underline{x}')\| \leq \kappa \|\underline{x} - \underline{x}'\|.$$

Then, the following hold true:

1. Let  $\underline{x}$  be a vector of  $\mathbb{R}^{(N-1)p}$  and  $\left(\underline{y}_{-j}\right)_{j \geq 1}$  a sequence of elements of  $E^q$ . Then the limit

$$\lim_{s \rightarrow \infty} G_{\underline{y}_1} \circ \dots \circ G_{\underline{y}_s}(\underline{x})$$

exists and does not depend on  $\underline{x}$ . Let  $H : (E^q)^\mathbb{N} \rightarrow \mathbb{R}^{(N-1)p}$  be the function defined by

$$H(\underline{y}_1, \underline{y}_2, \dots) = \lim_{s \rightarrow \infty} G_{\underline{y}_1} \circ \dots \circ G_{\underline{y}_s}(\underline{x}).$$

Then the function  $H$  is bounded. Moreover there exist  $C > 0$  such that

$$\|H(\underline{y}_1, \underline{y}_2, \dots) - H(\underline{y}'_1, \underline{y}'_2, \dots)\| \leq C \sum_{j \geq 1} \kappa^{j/k} \mathbb{1}_{\underline{y}_j \neq \underline{y}'_j}.$$

2. A process  $((Y_t, \lambda_t))_{t \in \mathbb{Z}}$  is solution of the problem  $\mathcal{P}_f$  is and only if  $(Y_t)_{t \in \mathbb{Z}}$  is a chain with complete connection associated to a function  $g$  (see Corollary 2) defined by

$$g(Y_{t-1}, Y_{t-2}, \dots) = H_1(V_t, V_{t-1}, \dots), \quad V_t = (Y_{t-1}, \dots, Y_{t-q}).$$

Here  $H_1$  denotes the  $N - 1$  first coordiantes of the function  $H$  defined previously.

3. There exists a unique strictly stationary solution to the equations (1) and (7). Moreover the process  $(Y_t)_{t \in \mathbb{Z}}$  is  $\phi$ -mixing with a geometric decrease for the mixing coefficients. This implies the ergodicity of the joint process  $((Y_t, \lambda_t))_{t \in \mathbb{Z}}$ .

*Proof.* 1. The first part of the assertion is a straightforward consequence of the assumption and is omitted. We focus on the proof of the Lipschitz property of the function  $H$ . For  $j \geq 1$ , we set

$$G_{\underline{y}}^{(j)} = G_{\underline{y}_{(j-1)k+1}} \circ G_{\underline{y}_{(j-1)k+2}} \circ \dots \circ G_{\underline{y}_{jk}}.$$

We have

$$H(\underline{y}_1, \underline{y}_2, \dots) = \lim_{s \rightarrow \infty} G_{\underline{y}}^{(1)} \circ \dots \circ G_{\underline{y}}^{(s)}(\underline{x}).$$

By the stated assumption, we obtain that

$$\|G_{\underline{y}}^{(j)}(x) - G_{\underline{y}'}^{(j)}(x)\| \leq \sum_{\ell=1}^k K^\ell \mathbb{1}_{\underline{y}_{(j-1)k+\ell} \neq \underline{y}'_{(j-1)k+\ell}}.$$

Hence

$$\|G_{\underline{y}}^{(1)} \circ \dots \circ G_{\underline{y}}^{(s)}(x) - G_{\underline{y}'}^{(1)} \circ \dots \circ G_{\underline{y}'}^{(s)}(x)\| \leq \sum_{j=1}^{\infty} K^j \sum_{\ell=1}^k K^\ell \mathbb{1}_{\underline{y}_{(j-1)k+\ell} \neq \underline{y}'_{(j-1)k+\ell}}.$$

By setting  $C = (K \vee 1)^k$  and letting  $x \rightarrow \infty$  we obtain the result.

2. From the first point of the theorem, the necessary condition follows easily. Now let us assume that  $\lambda_t = H_1(V_t, V_{t-1}, \dots)$ . Setting  $\underline{\lambda}_t = H(V_t, V_{t-1}, \dots)$ , note that the continuity of the function  $G_{V_t}$  implies that

$$\begin{aligned} \underline{\lambda}_t &= \lim_{s \rightarrow \infty} G_{V_t} \circ G_{V_{t-1}} \circ \dots \circ G_{V_{t-s}}(x) \\ &= G_{V_t} \left( \lim_{s \rightarrow \infty} G_{V_{t-1}} \circ \dots \circ G_{V_{t-s}}(x) \right) \\ &= G_{V_t}(\underline{\lambda}_t), \end{aligned}$$

which proofs that  $(\lambda_t)_{t \in \mathbb{Z}}$  satisfies (7).

3. The third point is a straightforward consequence of the two first results and of Corollary 2. Moreover the geometric decay of the  $\phi$ -mixing coefficients has been discussed in the remarks made following Proposition 2. Finally, it is well-known that  $\phi$ -mixing implies ergodicity of the process  $(Y_t)_{t \in \mathbb{Z}}$  and then ergodicity of the process  $((\lambda_t, Y_t))_{t \in \mathbb{Z}}$ ; see Samorodnitsky (2016, Ch. 2), for instance.

□

### 3.2 Linear models

Let  $A_0, A_1, \dots, A_p, B_1, \dots, B_q$  be some real matrices of size  $(N-1) \times (N-1)$ . We assume that

$$f(x_1, \dots, x_p; y_1, \dots, y_q) = A_0 + \sum_{i=1}^p A_i x_i + \sum_{i=1}^q B_i y_i.$$

Then the above model can be written alternatively as

$$G_{V_t}(\underline{x}) = \widetilde{A}\underline{x} + B,$$

with

$$\tilde{A} = \begin{pmatrix} A_1 & \cdots & A_{p-1} & A_p \\ & I_{(N-1)(p-1)} & & 0 \end{pmatrix}, \quad B = \begin{pmatrix} A_0 + \sum_{i=1}^q B_i Y_{t-i} \\ 0 \\ \vdots \\ 0 \end{pmatrix},$$

where  $I_{(N-1)(p-1)}$  denotes the identity matrix of order  $(N-1)(p-1)$ . Then, the assumptions of Theorem 1 are satisfied if the spectral radius of  $\tilde{A}$  is less than unity (and then the norm of  $\tilde{A}^k$  is less than one if  $k$  is large enough) which also means that the roots of the polynomial  $\mathcal{P}(z) = \det(I_{N-1} - \sum_{i=1}^p A_i z^i)$  are outside the unit disc; Lütkepohl (2005, Ch.2). For the case  $p = q = 1$ , this result improves the conditions proved by Moysiadis and Fokianos (2014) since it does not require an additional assumption on the coefficient  $B_1$ . In addition, the results answer in affirmative the question posed by Tjøstheim (2012) for the case of binary autoregressive model. Compared with the work of Fokianos and Moysiadis (2017) we note again that for the case of logistic autoregressive modeling with binary data, the obtained conditions simplify considerably since it is only required that the coefficients corresponding to the latent process have sum less than one in absolute value.

### 3.3 Non-linear models

Recall (7) and assume that there exists a norm  $\|\cdot\|$  on  $\mathbb{R}^{N-1}$  and some positive real numbers  $\beta_1, \dots, \beta_p, \alpha_1, \dots, \alpha_q$  with  $\alpha = \sum_{i=1}^p \alpha_i < 1$  and for all  $x_1, x'_1, \dots, x_p, x'_p, y_1, y'_1, \dots, y_q, y'_q \in \mathbb{R}^{N-1}$ ,

$$\|f(x_1, \dots, x_p; y_1, \dots, y_q) - f(x'_1, \dots, x'_p; y'_1, \dots, y'_q)\| \leq \sum_{i=1}^p \alpha_i \|x_i - x'_i\| + \sum_{i=1}^q \beta_i \|y_i - y'_i\|.$$

It can be proved under this condition that, for a process  $(Z_t)_{t \in \mathbb{Z}}$  taking values in  $E^q$ , the random mapping

$$G_{Z_t}(x) = (f(x, Z_t), x_1, \dots, x_{p-1})$$

is contracting, after iteration. Indeed, if  $x, y \in \mathbb{R}^p$ ,  $\sigma_i^x = x_i$  for  $1 \leq i \leq p$  and

$$\sigma_t^x = f(\sigma_{t-1}^x, \dots, \sigma_{t-p}^x, Z_t), \quad t \geq p+1,$$

it follows by induction that

$$\|\sigma_t^x - \sigma_t^y\| \leq \alpha^{\frac{t-p}{p}} \|x - y\|, \quad t \geq p+1.$$

Hence, there exists an integer  $m \geq 1$ , such that the mapping

$$H_t^{(m)}(x) = G_{Z_t} \circ G_{Z_{t-1}} \circ \cdots \circ G_{Z_{t-m}}(x)$$

satisfies

$$\left\| H_t^{(m)}(x) - H_t^{(m)}(x') \right\| \leq \kappa \|x - x'\|$$

for some  $\kappa \in (0, 1)$ . Therefore the assumption of Theorem 1 is satisfied. We note again that this condition improves upon the conditions obtained by Moysiadis and Fokianos (2014) and Fokianos and Moysiadis (2017) since they require only that  $\alpha < 1$ .

## 4 Inclusion of exogenous covariates

In this section, we study the problem of including a covariate process  $(Z_t)_{t \in \mathbb{Z}}$  in an autoregressive categorical time series model. We will be assuming that the covariate process is strongly exogenous. Such an assumption implies that for each time  $t$ ,  $Y_t$  is independent from  $(Z_s)_{s \geq t+1}$  conditionally to  $Y_{t-1}, Z_t, Y_{t-2}, Z_{t-1} \dots$  and allows simple computation of the likelihood function. Indeed, if  $f(\cdot | Y_{t-1}^-, Z_t)$  denotes the conditional density of  $Y_t$  given  $Y_{t-1}, Z_t, Y_{t-2}, Z_{t-1} \dots$ , then, for a bounded measurable function  $h$ ,

$$\mathbb{E} \left[ h(Y_1, Y_2, \dots, Y_n) \mid Z_1 = z_1, \dots, Z_n = z_n, Y_0^- = y_0^- \right] = \int h(y_1, \dots, y_n) \prod_{i=1}^n f(y_i | y_{i-1}^-, z_i) \mu(dy_1) \cdots \mu(dy_n),$$

where  $\mu$  denotes a reference measure for the model. This type of exogeneity is also called Granger causality or Sims causality in the literature (see, for instance, Gouriéroux and Monfort (1995, Sec. 1.5.2), for a discussion of these different concepts). We will be restricting our study to the case of finite order Markov chains, i.e. the parameter  $\lambda_t$  does not depend on its past values. The general case appears more difficult to tackle and is will be considered in another communication.

Note that even for a finite-state Markov chain, it is difficult to find in the literature a result which guarantees ergodicity when some covariates are included in the dynamic. As we will see, there is an interesting parallel between Markov chains with exogeneous covariates and Markov chains in random environments which were studied in probability theory. In the proof of Theorem 2 given below, we will use an approach discussed in Cogburn (1984) for showing ergodicity of Markov processes in random environments.

### 4.1 A general result for finite state Markov chains with covariates

We will be discussing results concerning stationarity and ergodicity of a finite state Markov chain which can be jointly observed with a covariate process. In what follows, denote by  $Z = (Z_t)_{t \in \mathbb{Z}}$  a stationary process with values in the space  $G = \mathbb{R}^d$  and  $(Y_t)_{t \in \mathbb{N}}$  a process which takes values in a finite set  $E$ . In addition,

conditionally on  $Z$ ,  $(Y_t)_{t \in \mathbb{N}}$  is a finite-state inhomogeneous Markov chain. More precisely, we assume that there exist a family of transition matrices  $\{P_g : g \in G\}$  such that

$$\mathbb{P}(Y_t = y | Y_{t-1} = x; Z) = P_{Z_t}(x, y), \quad (x, y) \in E^2. \quad (8)$$

Throughout the section we will assume the following:

**(E1)** There exists an integer  $m \geq 1$  such that for all  $(z_1, z_2, \dots, z_m, x, y) \in G^m \times E^2$ ,

$$P_{z_1} P_{z_2} \cdots P_{z_m}(x, y) > 0.$$

**(E2)** The process  $Z$  is mixing in the ergodic theory sense, i.e for all elements  $A$  and  $B$  of  $\mathcal{B}(G^{\mathbb{Z}})$ , we have

$$\lim_{n \rightarrow \infty} \mathbb{P}(Z \in A, \tau^n Z \in B) = \mathbb{P}(Z \in A) \mathbb{P}(Z \in B),$$

where  $\tau$  denotes the shift operator on  $F^{\mathbb{Z}}$  defined by  $\tau Z = (Z_{j+1})_{j \in \mathbb{Z}}$ .

Note that Assumption **(E1)** implies that a process  $(Y_t)_{t \in \mathbb{Z}}$  satisfying (8) also satisfies

$$\mathbb{P}(Y_{t+m} = y | Y_t = x, Z) > 0, \text{ a.s. } (x, y, t) \in E \times E \times \mathbb{Z}.$$

In addition, Assumption **(E2)** is stronger than the assuming ergodicity of the process  $Z$  but weaker than the classical strong mixing condition usually employed in the literature. A large number of useful stochastic processes are mixing, for instance the strong mixing processes and Bernoulli shifts defined by  $Z_t = H(\varepsilon_t, \varepsilon_{t-1}, \dots)$  where  $H$  is a measurable function and  $(\varepsilon_t)_{t \in \mathbb{Z}}$  is an i.i.d sequence. For instance, [Samorodnitsky \(2016, Ch.2\)](#), discusses several properties of the different types of mixing in ergodic theory of stationary processes. Assumption **(E2)** will be employed for obtaining ergodicity for the shift operator  $\tau^m$  which is not implied by the ergodicity of the shift operator  $\tau$ . The main result of this section is given by the following theorem.

**Theorem 2.** *Suppose that **(E1-E2)** hold true. Then there exists a unique stochastic process  $(Y_t)_{t \in \mathbb{Z}}$  satisfying (8). Moreover the process  $((Y_t, Z_t))_{t \in \mathbb{Z}}$  is ergodic.*

*Proof.* We first show that the almost sure limit  $\lim_{s \rightarrow \infty} P_{Z_{t-s}} \cdots P_{Z_t}(x, y)$  exists for each  $y \in E$  and does not depend on  $x$ . For Markov chains, this condition is comparable to the weak ergodicity notion, but here the limit is taken in the backward sense. See [Seneta \(2006\)](#) for several sufficient conditions ensuring weak

ergodicity properties of time-inhomogeneous Markov chains, using ergodicity coefficients. Recall that the so-called Dobrushin's contraction coefficient of a stochastic matrix  $P$  is defined by

$$c(P) = \frac{1}{2} \sup_{x \neq y \in E} \|P(x, \cdot) - P(y, \cdot)\|_{TV}.$$

We remind that for two probability measures  $\mu$  and  $\nu$  on the finite set  $E$ , the total variation distance between  $\mu$  and  $\nu$  is defined by  $\|\mu - \nu\|_{TV} = \sum_{x \in E} |\mu(x) - \nu(x)|$ . It is well known that we have the contraction

$$\|\mu P - \nu P\|_{TV} \leq c(P) \|\mu - \nu\|_{TV}$$

and for two stochastic matrices  $P$  and  $Q$ , we have  $c(PQ) \leq c(P)c(Q)$ .

Moreover  $c(P) \leq 1 - |E| \min_{x, y \in E} P(x, y)$ , where  $|E|$  denotes the cardinality of the set  $E$ . So Assumption **(E1)** ensures that for all  $t \in \mathbb{Z}$ ,  $c(P_{Z_t} P_{Z_{t+1}} \cdots P_{Z_{t+m-1}}) < 1$  a.s. Now let  $x \neq x' \in E$ ,  $t \in \mathbb{Z}$  and  $s = km + \ell$ . we obtain by setting  $\rho = 1 - \eta|E|$ ,

$$\|P_{Z_{t-s+1}} \cdots P_{Z_t}(x, \cdot) - P_{Z_{t-s+1}} \cdots P_{Z_t}(x', \cdot)\|_{TV} \leq 2c(P_{Z_{t-km+1}} \cdots P_{Z_t}) \leq 2 \prod_{j=0}^{k-1} c(P_{Z_{t-(j+1)m+1}} \cdots P_{Z_{t-jm}}).$$

From Assumption **(E2)**, the covariate process  $Z$  is mixing. Then the process  $(Z_{t-j})_{j \in \mathbb{Z}}$  is also mixing. Indeed, if  $\theta_t$  and  $\tau$  denote the applications defined on  $G^{\mathbb{Z}}$  by  $\theta_t x = (x_{t-i})_{i \in \mathbb{Z}}$  and  $\tau x = (x_{i+1})_{i \in \mathbb{Z}}$  respectively, we have  $\tau \circ \theta_t = \theta_t \circ \tau^{-1}$ . Then for two borelians  $A$  and  $B$ , we get

$$\begin{aligned} \mathbb{P}(\theta_t Z \in A, \tau^n \theta_t Z \in B) &= \mathbb{P}(Z \in \theta_t^{-1} A, \tau^{-n} Z \in \theta_t^{-1} B) \\ &= \mathbb{P}(\tau^n Z \in \theta_t^{-1} A, Z \in \theta_t^{-1} B) \\ &\rightarrow \mathbb{P}(\theta_t Z \in A) \mathbb{P}(\theta_t Z \in B). \end{aligned}$$

Moreover, observe that the operator  $\tau^m$  is ergodic for  $\mathbb{P}_Z$ . Indeed, if a borelian set  $A$  is such that  $\tau^m A = A$ , we have using assumption **E2**,

$$\mathbb{P}(Z \in A) = \mathbb{P}(\tau^{km} Z \in A, Z \in A) \rightarrow_{k \rightarrow \infty} \mathbb{P}(Z \in A)^2.$$

Then, we conclude that  $\mathbb{P}_Z(A) \in \{0, 1\}$ , which shows that  $\tau^m$  is ergodic. Now, using Assumption **(E1)**, we have  $\mathbb{E} \log c(P_{Z_1} \cdots P_{Z_m}) < 0$ . Then from the ergodic theorem, we get

$$\prod_{j=0}^{k-1} c(P_{Z_{t-(j+1)m+1}} \cdots P_{Z_{t-jm}}) = \exp \left( \sum_{j=0}^{k-1} \log c(P_{Z_{t-(j+1)m+1}} \cdots P_{Z_{t-jm}}) \right) \rightarrow_{k \rightarrow \infty} 0.$$

In addition, when  $n \geq s$ , we deduce that

$$\|P_{Z_{t-s+1}} \cdots P_{Z_t}(x, \cdot) - P_{Z_{t-n}} \cdots P_{Z_t}(x, \cdot)\|_{TV} \leq 2c(P_{Z_{t-s+1}} \cdots P_{Z_t}).$$

This shows that the product of matrices  $P_{Z_{t-s+1}} \cdots P_{Z_t}$  converges, when  $s \rightarrow \infty$ , to a stochastic matrix whose rows are all equal. Then there exists a measurable function  $D : F^{\mathbb{N}} \rightarrow E^N$  with  $N = |E|$  such that

$$D(Z_t, Z_{t-1}, \dots) = \lim_{s \rightarrow \infty} P_{Z_{t-s+1}} \cdots P_{Z_t}(x, \cdot) \text{ a.s.}$$

Setting  $\bar{D}_t = D(Z_t, Z_{t-1}, \dots)$ ,  $\bar{D}_t$  is a random probability measure on  $E$ . For  $t \in \mathbb{Z}$ ,  $z \in F^{\mathbb{Z}}$ ,  $k$  a non-negative integer and  $y_0, y_1, \dots, y_k \in E$ , we set

$$\mu_{t:t+k}(z; y_0, y_1, \dots, y_k) = \prod_{i=1}^k P_{z_{t+i}}(y_{i-1}, y_i) \bar{D}_t(y_0).$$

From the Kolmogorov extension theorem, there exists for  $\mathbb{P}_Z$ -almost all values of  $z \in \mathcal{B}(G^{\mathbb{Z}})$  a unique measure  $\mu(z, \cdot)$  on  $E^{\mathbb{Z}}$  with marginals  $\mu_{t:t+k}(z, \cdot)$ . Hence, if  $\zeta$  denotes the probability distribution of  $Z$ , the measure  $\gamma$  defined by

$$\gamma(A \times B) = \int_B \mu(z, A) \zeta(dz), \quad (A, B) \in \mathcal{B}(E^{\mathbb{Z}}) \times \mathcal{B}(F^{\mathbb{Z}}),$$

is that of a couple of stationary process  $(Y, Z)$  satisfying (8). To show uniqueness, let  $(Y'_t)_{t \in \mathbb{Z}}$  be another stochastic process satisfying (8). Then the distribution of  $Y'|Z = z$  is that of a non-homogeneous Markov chain with transitions  $\{P_{z_t} : t \in \mathbb{Z}\}$ . As shown before, this conditional distribution is unique and equal to  $\mu(z, \cdot)$ .

Next we show ergodicity of the process  $((Y_t, Z_t))_{t \in \mathbb{Z}}$ . To this end, we use an approach introduced in [Cogburn \(1984\)](#) for the study of Markov processes in random environment. This type of argument is also used in [Sinn and Poupart \(2011\)](#) for positive transition matrix  $P_{Z_t}$  and we give here a more general and shorter proof. The approach used in [Cogburn \(1984\)](#) consists in considering the Markov kernel  $Q$  on  $E \times F$  defined by

$$Q((x, z), \{y\} \times A) = P_{z_1}(x, y) \mathbb{1}_A(\tau z), \quad A \in \mathcal{B}(F^{\mathbb{Z}}).$$

If  $\nu$  denotes the probability distribution  $(Y_t, \tau^t Z)$  which takes values in  $E \times F$ , then  $\nu$  is invariant for  $Q$  and the process  $(H_t)_{t \in \mathbb{Z}}$  defined by  $H_t = (Y_t, \tau^t Z)$  is a Markov chain of transition kernel  $Q$ . Let  $C$  be a  $\nu$ -invariant set, i.e  $Q((x, z), C) = 1$  for  $\nu$ -almost every  $(x, z) \in C$ . Using Corollary 5.11 in [Hairer \(2006\)](#), the Markov chain  $(H_t)_{t \in \mathbb{Z}}$  forms an ergodic process if and only if every  $\nu$ -invariant set  $C$  is of measure 0 or 1. In our



case, we have  $C = \cup_{x \in E} \{x\} \times C_x$  for some  $C_x \in \mathcal{B}(F^{\mathbb{Z}})$ . To this end, we first note that if  $C$  is  $\nu$ -invariant, then

$$\nu(C) = \nu Q(C) = \int_C d\nu(x, z) Q((x, z), C) + \int_{C^c} d\nu(x, z) Q((x, z), C) = \nu(C) + \int_{C^c} d\nu(x, z) Q((x, z), C).$$

Then we get  $Q((x, z), C) = 0$  for  $\nu$ -almost every  $(x, z) \in C^c$ , the complement of  $C$  in  $E \times F$ . Hence, we obtain  $Q((x, z), C) = Q \mathbb{1}_C(x, z) = \mathbb{1}_C(x, z)$  for  $\nu$ -almost every  $(x, z)$ . But this also gives  $Q^m \mathbb{1}_C = \mathbb{1}_C$ ,  $\nu$  a.e, where  $m$  has been defined by assumption **(E1)**. Moreover, we have that

$$Q^m((x, z), C) = \sum_{y \in E} \mathbb{1}_C(y, \tau^m z) [P_{z_1} \cdots P_{z_m}](x, y). \quad (9)$$

We write  $A = B$   $\nu$ -a.e. if  $\nu(A \Delta B) = 0$  where  $A \Delta B$  denotes the symmetric difference of the sets  $A$  and  $B$ , i.e  $A \Delta B = (A \cap B^c) \cup (A^c \cap B)$ . From assumption **(E1)**, all the entries of the matrix  $P_{z_1} \cdots P_{z_m}$  are positive. Then we deduce that for almost every  $(x, z) \in C$ , we have  $(y, \tau^m z) \in C$  for all  $y \in E$ . We set  $D = \cap_{y \in E} C_y$ . Let us denote by  $\nu_1$  and  $\nu_2$  the marginals of  $\nu$ . We first remark that for all  $A \in \mathcal{B}(F)$ , we have

$$\nu(\{x\} \times A) = \int_A \mathbb{P}(Y_0 = x | Z = z) \nu_2(dz) = \sum_{y \in E} \int_A \mathbb{P}(Y_0 = x | Y_{-m} = y, Z = z) \mathbb{P}(Y_n = y | Z = z) \nu_2(dz).$$

Employing again assumption **(E1)**, we get  $\nu(\{x\} \times A) \geq \eta \nu_2(A)$ . For  $x \in E$ , we set  $B_x = \tau^m C_x \setminus D$ . As stated above we have

$$\nu(\{(x, z) : z \in B_x\}) = \sum_{x \in E} \nu(\{x\} \times B_x) = 0.$$

We conclude that  $\nu_2(B_x) = 0$  for all  $x \in E$  and then  $\nu_2(\tau^m C_x \setminus C_x) = 0$ . By stationarity,  $\nu_2(\tau^m C_x) = \nu_2(C_x)$ . Therefore, for every  $x \in E$ ,  $\tau^m C_x = C_x$ ,  $\mu$ -a.e. But using assumption **E2**, we have that

$$\nu_2(C_x) = \nu_2(\tau^{km} C_x \cap C_x) \rightarrow_{k \rightarrow \infty} \nu_2(C_x)^2.$$

Then, we conclude that  $\nu_2(C_x) \in \{0, 1\}$ . If  $\nu_2(C_x) = 0$  for every  $x$ , we easily get  $\nu(C) = 0$ . Now if there exists  $x \in E$  such that  $\nu_2(C_x) = 1$ , we have, using the equality  $\nu_2(B_x) = 0$ ,

$$1 \leq \nu_2(C_x) = \nu_2(\tau^m C_x) = \nu_2(\tau^m C_x \cap D) \leq \nu_2(D) \leq \min_{y \in E} \nu_2(C_y).$$

Then  $\nu_2(C_y) = 1$  for each  $y \in E$ . Finally we obtain

$$\nu(C) = \sum_{y \in E} \nu(\{y\} \times C_y) = \sum_{y \in E} \nu_1(y) = 1.$$

Hence, we have shown that the process  $(H_t)_{t \in \mathbb{Z}}$  is ergodic and so is the process  $((Y_t, Z_t)_{t \in \mathbb{Z}})$ .

□

## 4.2 Application to the multinomial logistic model with covariates

We assume here that conditionally to a covariate process  $(Z_t)_{t \in \mathbb{Z}}$  taking values in  $\mathbb{R}^d$ , the process  $(Y_t)_{t \in \mathbb{Z}}$  is a  $q$ -order Markov chain such that

$$\mathbb{P}(Y_t = e_j | Y_{t-1}, \dots, Y_{t-q}, Z) = \frac{\exp(g_j(Y_{t-1}, \dots, Y_{t-q}; Z_t))}{1 + \sum_{s=1}^{N-1} \exp(g_s(Y_{t-1}, \dots, Y_{t-q}; Z_t))} := Q_{Z_t}(Y_{t-q:t-1}, e_j)$$

for some measurable functions  $g_j : E^q \times \mathbb{R}^d$ ,  $1 \leq j \leq N-1$ . Let us check that assumption **E1** is satisfied for the conditional Markov chain  $(X_t)_{t \in \mathbb{Z}}$  defined by  $X_t = (Y'_t, Y'_{t-1}, \dots, Y'_{t-q+1})'$ . Conditionally to  $Z$ , the process  $(X_t)_{t \in \mathbb{Z}}$  defines a time-inhomogeneous Markov chains such that

$$P_{Z_t}((u_1, \dots, u_q), (v_1, \dots, v_q)) : = Q_{Z_t}((u_1, \dots, u_q), v_1) \prod_{s=1}^{q-1} \mathbb{1}_{v_{s+1}=u_s}.$$

Since the transition  $Q_{Z_t}$  takes only positive values, the assumption **E1** follows by taking  $m = q$  and  $\gamma = a^q$ . Then assuming **E2** for the covariate process, Theorem 2 applies and guarantees the ergodicity of the process  $((Y_t, Z_t))_{t \in \mathbb{Z}}$ .

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