# Detection of cellular aging in a Galton-Watson process 

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#### Abstract

We consider the bifurcating Markov chain model introduced by Guyon to detect cellular aging from cell lineage. To take into account the possibility for a cell to die, we use an underlying super-critical binary Galton-Watson process to describe the evolution of the cell lineage. We give in this more general framework a weak law of large number, an invariance principle and thus fluctuation results for the average over all individuals in a given generation, or up to a given generation. We also prove that the fluctuations over each generation are independent. Then we present the natural modifications of the tests given by Guyon in cellular aging detection within the particular case of the auto-regressive model.


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## 1. Introduction

This work is motivated by experiments done by biologists on Escherichia coli, see Stewart et al. [20]. E. coli is a rod-shaped single celled organism which reproduces by dividing in the middle. It produces a new end per progeny cell. We shall call this new end the new pole whereas the other end will be called the old pole. The age of a cell is given by the age of its old pole (i.e. the number of generations in the past of the cell before the old pole was produced). Notice that at each generation a cell gives birth to 2 cells which have a new pole and one of the two cells has an old pole of age one (which corresponds to the new pole of its mother), while the other has an old pole with age larger than one (which corresponds to the old pole of its mother). The

[^0]former is called the new pole daughter and the latter the old pole daughter. Experimental data, see [20], suggest strongly that the growth rate of the new pole daughter is significantly larger than the growth rate of the old pole daughter. For asymmetric aging see also [2] for an other case of asymmetric division, and Lindner et al. [15] or Ackermann et al. [1] on asymmetric damage repartition.

Guyon [11] studied a mathematical model, called bifurcating Markov chain (BMC), of an asymmetric Markov chain on a regular binary tree. This model allows to represent an asymmetric repartition for example of the growth rate of a cell between new pole and old pole daughters. Using this model, Guyon provides tests to detect a difference of the growth rate between new pole and old pole on a single experimental data set, whereas in [20] averages over many experimental data sets have to be done to detect this difference. In the BMC model, cells are assumed to never die (a death corresponds to no more division). Indeed few deaths appear in normal nutriment saturated conditions. However, under stress conditions, dead cells can represent a significant part of the population. It is therefore natural to take this random effect into account by using a Galton-Watson (GW) process. Our purpose is to study asymptotic results for bifurcating Markov chains on a Galton-Watson tree instead of a regular tree; see our main results in Section 1.4 (and also Theorems 3.7 and 5.2 for a more general model). Notice that inferences on symmetric bifurcating processes on regular trees have been studied, see the survey of Hwang, et al. [14] and the seminal work of Cowan and Staudte [9]. We also learned of a recent independent work on inferences for asymmetric auto-regressive models by Bercu, et al. [8]. Other models on cell lineage with differentiation have been investigated, see for example Bansaye [5,6] on parasite infection and Evans and Steinsaltz [10] on asymptotic models relying on super-Brownian motion. Besides, Markov chains on tree (random or not) have been widely studied. We mention the results of Yang [21] and Huang and Yang [13] on strong law of large numbers for at most countable Markov chains on non-random tree (homogeneous or not). See also [18] for a survey (in particular Section 3) on tree indexed processes. Athreya and Kang [3] have studied law of large numbers and its convergence rate for Markov chains on Galton-Watson tree, but the reproduction law does not charge 0 (each parent has at least one child). In those latter models, conditionally on the parent, the children behave independently: that is the division is symmetric. However, the correlation and the asymmetry between children, given the parent, is of main interest here. See $[2,15,1]$ for biological experiments. The present paper is concerned with law of large numbers and invariance principle for Markov chains indexed by a super-critical binary GW tree. We give in a paper with Bansaye and Tran [7] a partial extension of the present results to a continuous time setting.

One could ask if similar results hold for other random binary trees, for example critical or sub-critical GW trees conditioned to non-extinction (notice that in these two cases, the shape of the random tree is different from the super-critical one; in particular the asymptotic behavior of the number of individuals in one generation does not increase geometrically).

### 1.1. The statistical model

In order to study the behavior of the growth rate of cells in [20], we set some notations: we index the genealogical tree by the regular binary tree $\mathbb{T}=\{\emptyset\} \cup \bigcup_{k \in \mathbb{N}^{*}}\{0,1\}^{k} ; \emptyset$ is the label of the founder of the population and if $i$ denotes a cell, let $i 0$ denote the new pole progeny cell, and $i 1$ the old pole progeny cell. The growth rate of cell $i$ is $X_{i}$. When the mother gives birth to two cells among which a unique one divides, we consider that the cell which does not divide, does not grow. We study the growth rate of each cell generation by generation, using a discrete time

Markov chain described by the following model, which is a very simple case of the more general model of BMC on GW tree developed in Section 1.2:

- With probability $p_{1,0}, i$ gives birth to two cells $i 0$ and $i 1$ which will both divide. The growth rates of the daughters $X_{i 0}$ and $X_{i 1}$ are then linked to the mother's one $X_{i}$ through the following auto-regressive equations

$$
\left\{\begin{array}{l}
X_{i 0}=\alpha_{0} X_{i}+\beta_{0}+\varepsilon_{i 0}  \tag{1}\\
X_{i 1}=\alpha_{1} X_{i}+\beta_{1}+\varepsilon_{i 1},
\end{array}\right.
$$

where $\alpha_{0}, \alpha_{1} \in(-1,1), \beta_{0}, \beta_{1} \in \mathbb{R}$ and $\left(\left(\varepsilon_{i 0}, \varepsilon_{i 1}\right), i \in \mathbb{T}\right)$ is a sequence of independent centered bi-variate Gaussian random variables, with covariance matrix

$$
\sigma^{2}\left(\begin{array}{ll}
1 & \rho \\
\rho & 1
\end{array}\right), \quad \sigma^{2}>0, \rho \in(-1,1)
$$

- With probability $p_{0}$, only the new pole $i 0$ divides. Its growth rate $X_{i 0}$ is linked to its mother's one $X_{i}$ through the relation

$$
\begin{equation*}
X_{i 0}=\alpha_{0}^{\prime} X_{i}+\beta_{0}^{\prime}+\varepsilon_{i 0}^{\prime} \tag{2}
\end{equation*}
$$

where $\alpha_{0}^{\prime} \in(-1,1), \beta_{0}^{\prime} \in \mathbb{R}$ and $\left(\varepsilon_{i 0}^{\prime}, i \in \mathbb{T}\right)$ is a sequence of independent centered Gaussian random variables with variance $\sigma_{0}^{2}>0$.

- With probability $p_{1}$, only the old pole $i 1$ divides. Its growth rate $X_{i 1}$ is linked to its mother's one through the relation

$$
\begin{equation*}
X_{i 1}=\alpha_{1}^{\prime} X_{i}+\beta_{1}^{\prime}+\varepsilon_{i 1}^{\prime} \tag{3}
\end{equation*}
$$

where $\alpha_{1}^{\prime} \in(-1,1), \beta_{1}^{\prime} \in \mathbb{R}$ and $\left(\varepsilon_{i 1}^{\prime}, i \in \mathbb{T}\right)$ is a sequence of independent centered Gaussian random variables with variance $\sigma_{1}^{2}>0$.

- With probability $1-p_{1,0}-p_{1}-p_{0}$, which is non-negative, $i$ gives birth to two cells which do not divide.
- The sequences $\left(\left(\varepsilon_{i 0}, \varepsilon_{i 1}\right), i \in \mathbb{T}\right),\left(\varepsilon_{i 0}^{\prime}, i \in \mathbb{T}\right)$ and $\left(\varepsilon_{i 1}^{\prime}, i \in \mathbb{T}\right)$ are independent.

In Section 6, we first compute the maximum likelihood estimator (MLE) of the parameter

$$
\begin{equation*}
\theta=\left(\alpha_{0}, \beta_{0}, \alpha_{1}, \beta_{1}, \alpha_{0}^{\prime}, \beta_{0}^{\prime}, \alpha_{1}^{\prime}, \beta_{1}^{\prime}, p_{1,0}, p_{0}, p_{1}\right) \tag{4}
\end{equation*}
$$

and of $\kappa=\left(\sigma, \rho, \sigma_{0}, \sigma_{1}\right)$. Then, we prove that they are consistent (strong consistency can be achieved, see Remark 6.2), and that the MLE of $\theta$ is asymptotically normal, see Proposition 6.3 and Remark 6.5. Notice that the MLE of ( $p_{1,0}, p_{0}, p_{1}$ ), which is computed only on the underlying GW tree, was already known, see for example [16]. Eventually, we build a test for aging detection, for instance the null hypothesis $\left\{\left(\alpha_{0}, \beta_{0}\right)=\left(\alpha_{1}, \beta_{1}\right)\right\}$ against its alternative $\left\{\left(\alpha_{0}, \beta_{0}\right) \neq\left(\alpha_{1}, \beta_{1}\right)\right\}$, see Proposition 6.8. It appears that, for those hypothesis, using the test statistic from [11] with incomplete data due to death cells instead of the test statistic from Proposition 6.8 is not conservative, see Remark 6.9.

To prove those results, we shall consider a more general framework of BMC which is described in Section 1.2. An important tool is the auxiliary Markov chain which is defined in Section 1.3. Finally, easy to read version of our main general results are given in Section 1.4.

### 1.2. The mathematical model of bifurcating Markov chain (BMC)

We first introduce some notations related to the regular binary tree. Recall that $\mathbb{N}^{*}=\mathbb{N} \backslash\{0\}$, and let $\mathbb{G}_{0}=\{\emptyset\}, \mathbb{G}_{k}=\{0,1\}^{k}$ for $k \in \mathbb{N}^{*}, \mathbb{T}_{r}=\bigcup_{0 \leq k \leq r} \mathbb{G}_{k}$. The new (resp. old) pole daughter
of a cell $i \in \mathbb{T}$ is denoted by $i 0$ (resp. $i 1$ ), and 0 (resp. 1) if $i=\emptyset$ is the initial cell or root of the tree. The set $\mathbb{G}_{k}$ corresponds to all possible cells in the $k$-th generation. We denote by $|i|$ the generation of $i\left(|i|=k\right.$ if and only if $\left.i \in \mathbb{G}_{k}\right)$.

For a cell $i \in \mathbb{T}$, let $X_{i}$ denote a quantity of interest (for example its growth rate). We assume that the quantity of interest of the daughters of a cell $i$, conditionally on the generations previous to $i$, depends only on $X_{i}$. This property is stated using the formalism of BMC. More precisely, let $(E, \mathcal{E})$ be a measurable space, $P$ a probability kernel on $E \times \mathcal{E}^{2}: P(\cdot, A)$ is measurable for all $A \in \mathcal{E}^{2}$, and $P(x, \cdot)$ is a probability measure on $\left(E^{2}, \mathcal{E}^{2}\right)$ for all $x \in E$. For any measurable real-valued bounded function $g$ defined on $E^{3}$ we set

$$
(P g)(x)=\int_{E^{2}} g(x, y, z) P(x, \mathrm{~d} y, \mathrm{~d} z) .
$$

When there is no possible confusion, we shall write $P g(x)$ for $(P g)(x)$ to simplify notations.
Definition 1.1. We say a stochastic process indexed by $\mathbb{T}, X=\left(X_{i}, i \in \mathbb{T}\right)$, is a bifurcating Markov chain on a measurable space $(E, \mathcal{E})$ with initial distribution $v$ and probability kernel $P$, a $P$-BMC in short, if:

- $X_{\emptyset}$ is distributed as $\nu$.
- For any measurable real-valued bounded functions $\left(g_{i}, i \in \mathbb{T}\right)$ defined on $E^{3}$, we have for all $k \geq 0$,

$$
\mathbb{E}\left[\prod_{i \in \mathbb{G}_{k}} g_{i}\left(X_{i}, X_{i 0}, X_{i 1}\right) \mid \sigma\left(X_{j} ; j \in \mathbb{T}_{k}\right)\right]=\prod_{i \in \mathbb{G}_{k}} P g_{i}\left(X_{i}\right) .
$$

We consider a metric measurable space $(S, \mathscr{S})$ and add a cemetery point to $S, \partial$. Let $\bar{S}=S \cup\{\partial\}$, and $\overline{\mathscr{S}}$ be the $\sigma$-field generated by $\mathscr{S}$ and $\{\partial\}$. (In the biological framework of the previous Section, $S$ corresponds to the state space of the quantity of interest, and $\partial$ is the default value for dead cells.) Let $P^{*}$ be a probability kernel defined on $\bar{S} \times \overline{\mathscr{S}}^{2}$ such that

$$
\begin{equation*}
P^{*}(\partial,\{(\partial, \partial)\})=1 . \tag{5}
\end{equation*}
$$

Notice that this condition means that $\partial$ is an absorbing state. (In the biological framework of the previous Section, condition (5) states that no dead cell can give birth to a living cell.)

Definition 1.2. Let $X=\left(X_{i}, i \in \mathbb{T}\right)$ be a $P^{*}$-BMC on ( $\left.\bar{S}, \overline{\mathscr{S}}\right)$, with $P^{*}$ satisfying (5). We call $\left(X_{i}, i \in \mathbb{T}^{*}\right)$, with $\mathbb{T}^{*}=\left\{i \in \mathbb{T}: X_{i} \neq \partial\right\}$, a bifurcating Markov chain on a Galton-Watson tree. The $P^{*}$-BMC is said spatially homogeneous if $p_{1,0}=P^{*}(x, S \times S), p_{0}=P^{*}(x, S \times\{\partial\})$ and $p_{1}=P^{*}(x,\{\partial\} \times S)$ do not depend on $x \in S$. A spatially homogeneous $P^{*}$-BMC is said super-critical if $m>1$, where $m=2 p_{1,0}+p_{1}+p_{0}$.

Notice that condition (5) and the spatial homogeneity property imply that $\mathbb{T}^{*}$ is a GW tree. This justifies the name of BMC on a Galton-Watson tree. The GW tree is super-critical if and only if $m>1$. From now on, we shall only consider super-critical spatially homogeneous $P^{*}$-BMC on a Galton-Watson tree. (In the biological framework of the previous Section, $\mathbb{T}^{*}$ denotes the sub-tree of living cells and the notations $p_{1,0}, p_{0}$ and $p_{1}$ are consistent since, for instance, $P^{*}(x, S \times S)$ represents the probability that a living cell with growing rate $x$ gives birth to two living cells.)

We now consider the Galton-Watson sub-tree $\mathbb{T}^{*}$. For any subset $J \subset \mathbb{T}$, let

$$
\begin{equation*}
J^{*}=J \cap \mathbb{T}^{*}=\left\{j \in J: X_{j} \neq \partial\right\} \tag{6}
\end{equation*}
$$

be the subset of living cells among $J$, and $|J|$ be the cardinal of $J$. The process $Z=\left(Z_{k}, k \in \mathbb{N}\right)$, where $Z_{k}=\left|\mathbb{G}_{k}^{*}\right|$, is a GW process with reproduction generating function

$$
\psi(z)=\left(1-p_{0}-p_{1}-p_{1,0}\right)+\left(p_{0}+p_{1}\right) z+p_{1,0} z^{2} .
$$

Notice the average number of daughters alive is $m$. We have, for $k \geq 0$,

$$
\begin{equation*}
\mathbb{E}\left[\left|\mathbb{G}_{k}^{*}\right|\right]=m^{k} \quad \text { and } \quad \mathbb{E}\left[\left|\mathbb{T}_{r}^{*}\right|\right]=\sum_{q=0}^{r} \mathbb{E}\left[\left|\mathbb{G}_{q}^{*}\right|\right]=\sum_{q=0}^{r} m^{q}=\frac{m^{r+1}-1}{m-1} \tag{7}
\end{equation*}
$$

Let us recall some well-known facts on super-critical GW, see e.g. [12] or [4]. The extinction probability of the GW process $Z$ is $\eta=\mathbb{P}\left(\left|\mathbb{T}^{*}\right|<\infty\right)=1-\frac{m-1}{p_{1,0}}$. There exists a non-negative random variable $W$ s.t.

$$
\begin{equation*}
W=\lim _{q \rightarrow \infty} m^{-q}\left|\mathbb{G}_{q}^{*}\right| \quad \text { a.s. and in } L^{2} \tag{8}
\end{equation*}
$$

$\mathbb{P}(W=0)=\eta$ and whose Laplace transform, $\varphi(\lambda)=\mathbb{E}\left[\mathrm{e}^{-\lambda W}\right]$, satisfies $\varphi(\lambda)=\psi(\varphi(\lambda / m))$ for $\lambda \geq 0$. Notice the distribution of $W$ is completely characterized by this functional equation and $\mathbb{E}[W]=1$.

For $i \in \mathbb{T}$, we set $\Delta_{i}=\left(X_{i}, X_{i 0}, X_{i 1}\right)$, the mother-daughters quantities of interest. For a finite subset $J \subset \mathbb{T}$, we set

$$
M_{J}(f)= \begin{cases}\sum_{i \in J} f\left(X_{i}\right) & \text { for } f \in \mathcal{B}(\bar{S})  \tag{9}\\ \sum_{i \in J} f\left(\Delta_{i}\right) & \text { for } f \in \mathcal{B}\left(\bar{S}^{3}\right)\end{cases}
$$

with the convention that $M_{\emptyset}(f)=0$. We also define the following two averages of $f$ over $J$

$$
\begin{equation*}
\bar{M}_{J}(f)=\frac{1}{|J|} M_{J}(f) \quad \text { if }|J|>0 \quad \text { and } \quad \tilde{M}_{J}(f)=\frac{1}{\mathbb{E}[|J|]} M_{J}(f) \quad \text { if } \mathbb{E}[|J|]>0 \tag{10}
\end{equation*}
$$

We shall study the limit of the averages of a function $f$ of the BMC over the $n$-th generation, $\bar{M}_{\mathbb{G}_{n}^{*}}(f)$ and $\widetilde{M}_{\mathbb{G}_{n}^{*}}(f)$, or over all the generations up to the $n$-th, $\bar{M}_{\mathbb{T}_{n}^{*}}(f)$ and $\widetilde{M}_{\mathbb{T}_{n}^{*}}(f)$, as $n$ goes to infinity. Notice the no death case studied in [11] corresponds to $p_{1,0}=1$, that is $m=2$.

### 1.3. The auxiliary Markov chain

We define the sub-probability kernel on $S \times \mathscr{S}^{2}: P(\cdot, \cdot)=P^{*}\left(\cdot, \cdot \bigcap S^{2}\right)$, and two subprobability kernels on $S \times \mathscr{S}$ :

$$
P_{0}^{*}=P^{*}(\cdot,(\cdot \bigcap S) \times \bar{S}) \quad \text { and } \quad P_{1}^{*}=P^{*}(\cdot, \bar{S} \times(\cdot \bigcap S))
$$

Notice that $P_{0}^{*}$ (resp. $P_{1}^{*}$ ) is the restriction of the first (resp. second) marginal of $P^{*}$ to $S$. From spatial homogeneity, we have for all $x \in S, P^{*}\left(x, S^{2}\right)=p_{1,0}$ and, for $\delta \in\{0,1\}$,

$$
P_{\delta}^{*}(x,\{\partial\})=0 \quad \text { and } \quad P_{\delta}^{*}(x, S)=p_{\delta}+p_{1,0} .
$$

We introduce an auxiliary Markov chain (see [11] for the case $m=2$ ). Let $Y=\left(Y_{n}, n \in \mathbb{N}\right)$ be a Markov chain on $S$ with $Y_{0}$ distributed as $X_{\emptyset}$ and transition kernel

$$
Q=\frac{1}{m}\left(P_{0}^{*}+P_{1}^{*}\right) .
$$

The distribution of $Y_{n}$ corresponds to the distribution of $X_{I}$ conditionally on $\left\{I \in \mathbb{T}^{*}\right\}$, where $I$ is chosen at random in $\mathbb{G}_{n}$, see Lemma 2.1 for a precise statement. We shall write $\mathbb{E}_{x}$ when $X_{\emptyset}=x$ (i.e. initial distribution $v$ is the Dirac mass at $x \in S$ ).

Last, we need some more notation: if $(E, \mathcal{E})$ is a metric measurable space, then $\mathcal{B}_{b}(E)$ (resp. $\mathcal{B}_{+}(E)$ ) denotes the set of bounded (resp. non-negative) real-valued measurable functions defined on $E$. The set $\mathcal{C}_{b}(E)$ (resp. $\left.\mathcal{C}_{+}(E)\right)$ denotes the set of bounded (resp. non-negative) real-valued continuous functions defined on $E$. For a finite measure $\lambda$ on $(E, \mathcal{E})$ and $f \in$ $\mathcal{B}_{b}(E) \cup \mathcal{B}_{+}(E)$ we shall write $\langle\lambda, f\rangle$ for $\int f(x) \mathrm{d} \lambda(x)$.

We consider the following hypothesis (H):
The Markov chain $Y$ is ergodic, that is there exists a probability measure $\mu$ on $(S, \mathscr{S})$ s.t., for all $f \in \mathcal{C}_{b}(S)$ and all $x \in S, \lim _{k \rightarrow \infty} \mathbb{E}_{x}\left[f\left(Y_{k}\right)\right]=\langle\mu, f\rangle$.

Notice that under (H), the probability measure $\mu$ is the unique stationary distribution of $Y$ and ( $Y_{n}, n \in \mathbb{N}$ ) converges in distribution to $\mu$.

Remark 1.3. If we assume that $\left(\left(X_{i 0}, X_{i 1}\right), i \in \mathbb{T}\right)$ is ergodic, in the sense that there exists a probability measure $v$ on $\left(S^{2}, \mathscr{S}^{2}\right)$ s.t., for all $g \in \mathcal{C}_{b}\left(S^{2}\right)$ and all $(y, z) \in S^{2}$,

$$
\lim _{|i| \rightarrow \infty} \mathbb{E}_{(y, z)}\left[g\left(X_{i 0}, X_{i 1}\right)\right]=\langle v, g\rangle,
$$

then $Y$ is also ergodic.

### 1.4. The main results

We can now state our principal results on the weak law of large numbers and fluctuations for the averages over a generation or up to a generation. Those results are a particular case of the more general statements given in Theorems 3.7 and 5.2, using Remark 2.2.

Theorem 1.4. Let $\left(X_{i}, i \in \mathbb{T}^{*}\right)$ be a super-critical spatially homogeneous $P^{*}$-BMC on a $G W$ tree and $W$ be defined by (8). We assume that $(H)$ holds and that $x \mapsto P^{*} g(x) \in \mathcal{C}_{b}(\bar{S})$ for all $g \in \mathcal{C}_{b}\left(\bar{S}^{3}\right)$. Let $f \in \mathcal{C}_{b}\left(\bar{S}^{3}\right)$.

- Weak law of large numbers. We have the following convergence in probability:

$$
\begin{align*}
& \mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}} \frac{1}{\left|\mathbb{G}_{r}^{*}\right|} \sum_{i \in \mathbb{G}_{r}^{*}} f\left(\Delta_{i}\right)=\mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}} \bar{M}_{\mathbb{G}_{r}^{*}}(f) \xrightarrow[r \rightarrow \infty]{\mathbb{P}}\left\langle\mu, P^{*} f\right\rangle \mathbf{1}_{\{W \neq 0\}},  \tag{11}\\
& \mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}} \frac{1}{\left|\mathbb{T}_{r}^{*}\right|} \sum_{i \in \mathbb{T}_{r}^{*}} f\left(\Delta_{i}\right)=\mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}} \bar{M}_{\mathbb{T}_{r}^{*}}(f) \xrightarrow[r \rightarrow \infty]{ }\left\langle\mu, P^{*} f\right\rangle \mathbf{1}_{\{W \neq 0\}} . \tag{12}
\end{align*}
$$

- Fluctuations. We have the following convergence in distribution:

$$
\mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}} \frac{1}{\sqrt{\left|\mathbb{T}_{r}^{*}\right|}} \sum_{i \in \mathbb{T}_{r}^{*}}\left(f\left(\Delta_{i}\right)-P^{*} f\left(X_{i}\right)\right) \xrightarrow[r \rightarrow \infty]{(d)} \mathbf{1}_{\{W \neq 0\}} \sigma G,
$$

where $\sigma^{2}=\left\langle\mu, P^{*}\left(f^{2}\right)-\left(P^{*} f\right)^{2}\right\rangle$, and $G$ is a Gaussian random variable with mean zero, variance 1 , and independent of $W$.

Remark 1.5. The weak laws of large numbers given by (11) and (12) can be turned into strong laws of large numbers under stronger hypothesis, see Theorem 3.8.

We also can prove that the fluctuations over each generation are asymptotically independent.
Theorem 1.6. Let $\left(X_{i}, i \in \mathbb{T}^{*}\right)$ be a super-critical spatially homogeneous $P^{*}$-BMC on a $G W$ tree. We assume that $(H)$ holds and that $x \mapsto P^{*} g(x) \in \mathcal{C}_{b}(\bar{S})$ for all $g \in \mathcal{C}_{b}\left(\bar{S}^{3}\right)$. Let $d \geqq 1$, and for $\ell \in\{1, \ldots, d\}, f_{\ell} \in \mathcal{C}_{b}\left(\bar{S}^{3}\right)$ and $\sigma_{\ell}^{2}=\left\langle\mu, P^{*}\left(f_{\ell}^{2}\right)-\left(P^{*} f_{\ell}\right)^{2}\right\rangle$. We set for $f \in \mathcal{C}_{b}\left(\bar{S}^{3}\right)$

$$
N_{n}(f)=\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}} \frac{1}{\sqrt{\left|\mathbb{G}_{n}^{*}\right|}} \sum_{i \in \mathbb{G}_{n}^{*}}\left(f\left(\Delta_{i}\right)-P^{*} f\left(X_{i}\right)\right) .
$$

Then we have the following convergence in distribution:

$$
\left(N_{n}\left(f_{1}\right), \ldots, N_{n-d+1}\left(f_{d}\right)\right) \xrightarrow[n \rightarrow \infty]{(d)} \mathbf{1}_{\{W \neq 0\}}\left(\sigma_{1} G_{1}, \ldots, \sigma_{d} G_{d}\right),
$$

where $G_{1}, \ldots, G_{d}$ are independent Gaussian random variables with mean zero and variance 1, and are independent of $W$ given by (8).

Even if the results on fluctuations in Theorem 1.4 are not as complete as one might hope (see Remark 1.7), they are still sufficient to study the statistical model we gave in Section 1.1 for the detection of cellular aging from cell lineage when death of cells can occur.

Remark 1.7. Let $V=\left(V_{r}, r \geq 0\right)$ be a Markov chain on a finite state space. We assume $V$ is irreducible, with transition matrix $R$ and unique invariant distribution $\mu$. Then it is well known, see [17], that $\frac{1}{r} \sum_{i=1}^{r} h\left(V_{i}\right)$ converges a.s. to $\langle\mu, h\rangle$ and that, to prove the fluctuations result, one solves the Poisson equation $H-R H=h-\langle\mu, h\rangle$, writes

$$
\begin{align*}
\frac{1}{\sqrt{r}} \sum_{i=1}^{r}\left(h\left(V_{i}\right)-\langle\mu, h\rangle\right)= & \frac{1}{\sqrt{r}} \sum_{i=1}^{r}\left(H\left(V_{i}\right)-R H\left(V_{i-1}\right)\right) \\
& +\frac{1}{\sqrt{r}} R H\left(V_{0}\right)-\frac{1}{\sqrt{r}} R H\left(V_{r}\right) \tag{13}
\end{align*}
$$

and then uses martingale theory to obtain the asymptotic normality of $\frac{1}{\sqrt{r}} \sum_{i=1}^{r} H\left(V_{i}\right)-$ $R H\left(V_{i-1}\right)$ (we use similar techniques to prove the fluctuations in Theorem 1.4). It then only remains to say that $\frac{1}{\sqrt{r}} R H\left(V_{0}\right)$ and $\frac{1}{\sqrt{r}} R H\left(V_{r}\right)$ converge to 0 to conclude.

Assume that hypothesis of Theorem 1.4 hold and that $x \mapsto P^{*}(x, A)$ is continuous for all $A \in \mathcal{B}\left(\bar{S}^{2}\right)$. Let $h \in \mathcal{C}_{b}(\bar{S})$. Theorem 1.4 implies that $\mathbf{1}_{\left\{\left|\left|\mathbb{G}_{r}^{*}\right|>0\right\}\right.} \frac{1}{\left|\mathbb{T}_{r \mid}^{*}\right|} \sum_{i \in \mathbb{T}_{r}^{*}} h\left(X_{i}\right)$ converges in probability to $\langle\mu, h\rangle \mathbf{1}_{\{W \neq 0\}}$. To get the fluctuations, that is the limit of

$$
\mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}} \frac{1}{\sqrt{\left|\mathbb{T}_{r}^{*}\right|}} \sum_{i \in \mathbb{T}_{r}^{*}}\left(h\left(X_{i}\right)-\langle\mu, h\rangle\right)
$$

as $r$ goes to infinity, using martingale theory, one can think of using the same kind of approach in order to use the result on fluctuations of Theorem 1.4. But then notice that what will correspond to the boundary term in (13) at time $r, \frac{1}{\sqrt{r}} R H\left(V_{r}\right)$, will now be a boundary term over the last generation $\mathbb{G}_{r}^{*}$, whose cardinal is of the same order as $\left|\mathbb{T}_{r}^{*}\right|$. Thus the order of the boundary term is not negligible, and we cannot conclude using this approach.

The fluctuations for $\sum_{i \in \mathbb{T}_{r}^{*}} h\left(X_{i}\right)$ are still an open question.

### 1.5. Organization of the paper

We quickly study the auxiliary chain in Section 2 . We state the results on the weak law of large number in Section 3. Section 4 is devoted to some preparatory results in order to apply results on fluctuations for martingale. Our main result, Theorem 5.2, is stated and proved in Section 5. The biological model of Section 1.1 is analysed in Section 6.

## 2. Preliminary result and notations

Recall the Markov chain $Y$ defined in Section 1.3.
Lemma 2.1. We have, for $f \in \mathcal{B}_{b}(S) \cup \mathcal{B}_{+}(S)$,

$$
\begin{align*}
\mathbb{E}\left[f\left(Y_{n}\right)\right] & =m^{-n} \sum_{i \in \mathbb{G}_{n}} \mathbb{E}\left[f\left(X_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right]=\frac{\sum_{i \in \mathbb{G}_{n}} \mathbb{E}\left[f\left(X_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right]}{\sum_{i \in \mathbb{G}_{n}} \mathbb{P}\left(i \in \mathbb{T}^{*}\right)} \\
& =\mathbb{E}\left[f\left(X_{I}\right) \mid I \in \mathbb{T}^{*}\right], \tag{14}
\end{align*}
$$

where I is a uniform random variable on $\mathbb{G}_{n}$ independent of $X$.
Proof. We consider the first equality. Recall that $Y_{0}$ has distribution $\nu$. For $i=i_{1} \ldots i_{n} \in \mathbb{G}_{n}$, we have, thanks to (5) and the definition of $P^{*}$,

$$
\mathbb{E}\left[f\left(X_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right]=\mathbb{E}\left[f\left(X_{i}\right) \mathbf{1}_{\left\{X_{i} \neq \partial\right\}}\right]=\left\langle\nu,\left(P_{i_{1}}^{*} \ldots P_{i_{n}}^{*}\right) f\right\rangle,
$$

so that

$$
\begin{aligned}
\sum_{i \in \mathbb{G}_{n}} \mathbb{E}\left[f\left(X_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right] & =\sum_{i_{1}, \ldots, i_{n} \in\{0,1\}}\left\langle v,\left(P_{i_{1}}^{*} \ldots P_{i_{n}}^{*}\right) f\right\rangle \\
& =\left\langle v,\left(P_{0}^{*}+P_{1}^{*}\right)^{n} f\right\rangle=m^{n}\left\langle v, Q^{n} f\right\rangle=m^{n} \mathbb{E}\left[f\left(Y_{n}\right)\right] .
\end{aligned}
$$

This gives the first equality. Then take $f=1$ in the previous equality to get $m^{n}=\sum_{i \in \mathbb{G}_{n}}$ $\mathbb{P}\left(i \in \mathbb{T}^{*}\right)$ and the second equality of (14). The last equality of (14) is obvious.

We recall that $v$ denotes the distribution of $X_{\emptyset}$. Any function $f$ defined on $S$ is extended to $\bar{S}$ by setting $f(\partial)=0$. Let $F$ be a vector subspace of $\mathcal{B}(S)$ s.t.
(i) $F$ contains the constants;
(ii) $F^{2}:=\left\{f^{2} ; f \in F\right\} \subset F$;
(iii) (a) $F \otimes F \subset L^{1}(P(x, \cdot))$ for all $x \in S$ and $P\left(f_{0} \otimes f_{1}\right) \in F$ for all $f_{0}, f_{1} \in F$;
(b) For $\delta \in\{0,1\}, F \subset L^{1}\left(P_{\delta}^{*}(x, \cdot)\right)$ for all $x \in S$ and $P_{\delta}^{*}(f) \in F$ for all $f \in F$;
(iv) There exists a probability measure $\mu$ on $(S, \mathscr{S})$ s.t. $F \subset L^{1}(\mu)$ and $\lim _{n \rightarrow \infty} \mathbb{E}_{x}\left[f\left(Y_{n}\right)\right]=$ $\langle\mu, f\rangle$ for all $x \in S$ and $f \in F$;
(v) For all $f \in F$, there exists $g \in F$ s.t. for all $r \in \mathbb{N},\left|Q^{r} f\right| \leq g$;
(vi) $F \subset L^{1}(v)$.

By convention a function defined on $\bar{S}$ is said to belong to $F$ if its restriction to $S$ belongs to $F$.

Remark 2.2. Notice that if $(H)$ is satisfied and if $x \mapsto P^{*} g(x)$ is continuous on $S$ for all $g \in \mathcal{C}_{b}\left(\bar{S}^{3}\right)$ then the set $\mathcal{C}_{b}(S)$ fulfills (i)-(vi).

## 3. Weak law of large numbers

We give the first result of this section. Recall notations (6), (9) and (10).
Theorem 3.1. Let $\left(X_{i}, i \in \mathbb{T}^{*}\right)$ be a super-critical spatially homogeneous $P^{*}$-BMC on a $G W$ tree. Let $F$ satisfy (i)-(vi) and $f \in F$. The sequence $\left(\widetilde{M}_{\mathbb{G}_{q}^{*}}(f), q \in \mathbb{N}\right)$ converges to $\langle\mu, f\rangle W$ in $L^{2}$ as $q \rightarrow \infty$, where $W$ is defined by (8). We also have that the sequence $\left(\bar{M}_{\mathbb{G}_{q}^{*}}(f) \mathbf{1}_{\left\{\left|\mathbb{G}_{q}^{*}\right|>0\right\}}\right.$, $q \in \mathbb{N}$ ) converges to $\langle\mu, f\rangle \mathbf{1}_{\{W \neq 0\}}$ in probability when $q \rightarrow \infty$.

Remark 3.2. Intuitively, one can understand the result as follows. For one individual $i$ picked at random in $\mathbb{G}_{q}^{*}$, we get that $X_{i}$ is distributed as $Y_{q}$, and thus, when $q$ is large, as its stationary distribution $\mu$. Furthermore, two individuals $i$ and $j$ picked at random in $\mathbb{G}_{q}^{*}$ have, with high probability, their most recent common ancestor in one of the first generations. This implies that $X_{i}$ and $X_{j}$ are almost independent, thanks to ergodic property. In conclusion, the average of $f\left(X_{i}\right)$ over $\mathbb{G}_{q}^{*}$ behaves like $\langle\mu, f\rangle$.

One can also get an a.s. convergence in Theorem 3.1 under stronger hypothesis on $Y$ (such as geometric ergodicity) using similar arguments as in [11], see Theorem 3.8.

Proof. We first assume that $\langle\mu, f\rangle=0$. We have,

$$
\begin{aligned}
\left\|\sum_{i \in \mathbb{G}_{q}^{*}} f\left(X_{i}\right)\right\|_{L^{2}}^{2} & =\mathbb{E}\left[\left(\sum_{i \in \mathbb{G}_{q}} f\left(X_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right)^{2}\right]=\sum_{i \in \mathbb{G}_{q}} \mathbb{E}\left[f^{2}\left(X_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right]+B_{q} \\
& =m^{q} \mathbb{E}\left[f^{2}\left(Y_{q}\right)\right]+B_{q}
\end{aligned}
$$

with $B_{q}=\sum_{(i, j) \in \mathbb{G}_{q}^{2}, i \neq j} \mathbb{E}\left[f\left(X_{i}\right) f\left(X_{j}\right) \mathbf{1}_{\left\{(i, j) \in \mathbb{T}^{* 2}\right\}}\right]$, where we used (14) for the last equality.
Since the sum in $B_{q}$ concerns all pairs of distinct elements of $\mathbb{G}_{q}$, we have that $i \wedge j$, the most recent common ancestor of $i$ and $j$, does not belong to $\mathbb{G}_{q}$. We shall compute $B_{q}$ by decomposing this sum according to the generation of $k=i \wedge j: B_{q}=\sum_{r=0}^{q-1} \sum_{k \in \mathbb{G}_{r}} C_{k}$ with

$$
C_{k}=\sum_{(i, j) \in \mathbb{G}_{q}^{2}, i \wedge j=k} \mathbb{E}\left[f\left(X_{i}\right) f\left(X_{j}\right) \mathbf{1}_{\left\{(i, j) \in \mathbb{T}^{* 2}\right\}}\right]
$$

If $|k|=q-1$, using the Markov property of $X$ and of the GW process at generation $q-1$, we get

$$
\begin{aligned}
C_{k} & =\sum_{(i, j) \in \mathbb{G}_{1}^{2}, i \wedge j=\emptyset} \mathbb{E}\left[\mathbb{E}_{X_{k}}\left[f\left(X_{i}\right) f\left(X_{j}\right) \mathbf{1}_{\left\{(i, j) \in \mathbb{T}^{* 2}\right\}}\right] \mathbf{1}_{\left\{k \in \mathbb{T}^{*}\right\}}\right] \\
& =2 \mathbb{E}\left[P(f \otimes f)\left(X_{k}\right) \mathbf{1}_{\left\{k \in \mathbb{T}^{*}\right\}}\right]
\end{aligned}
$$

If $|k|<q-1$, we have, with $r=|k|$,

$$
\begin{aligned}
C_{k} & =2 \sum_{(i, j) \in \mathbb{G}_{q-r-1}^{2}} \mathbb{E}\left[\mathbb{E}_{X_{k 0}}\left[f\left(X_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right] \mathbb{E}_{X_{k 1}}\left[f\left(X_{j}\right) \mathbf{1}_{\left\{j \in \mathbb{T}^{*}\right\}}\right] \mathbf{1}_{\left\{k 0 \in \mathbb{T}^{*}, k 1 \in \mathbb{T}^{*}\right\}}\right] \\
& =2 \mathbb{E}\left[\sum_{i \in \mathbb{G}_{q-r-1}} \mathbb{E}_{X_{k 0}}\left[f\left(X_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right] \sum_{j \in \mathbb{G}_{q-r-1}} \mathbb{E}_{X_{k 1}}\left[f\left(X_{j}\right) \mathbf{1}_{\left\{j \in \mathbb{T}^{*}\right\}}\right] \mathbf{1}_{\left\{k 0 \in \mathbb{T}^{*}, k 1 \in \mathbb{T}^{*}\right\}}\right] \\
& =2 m^{2(q-r-1)} \mathbb{E}\left[\mathbb{E}_{X_{k 0}}\left[f\left(Y_{q-r-1}\right)\right] \mathbb{E}_{X_{k 1}}\left[f\left(Y_{q-r-1}\right)\right] \mathbf{1}_{\left\{k 0 \in \mathbb{T}^{*}, k 1 \in \mathbb{T}^{*}\right\}}\right] \\
& =2 m^{2(q-r-1)} \mathbb{E}\left[P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)\left(X_{k}\right) \mathbf{1}_{\left\{k \in \mathbb{T}^{*}\right\}}\right],
\end{aligned}
$$

where we used the Markov property of $X$ and of the GW process at generation $r+1$ for the first equality, (14) for the third equality, and the Markov property at generation $r$ for the last equality.

In particular, we get that $C_{k}=2 m^{2(q-r-1)} \mathbb{E}\left[P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)\left(X_{k}\right) \mathbf{1}_{\left\{k \in \mathbb{T}^{*}\right\}}\right]$ for all $k$ s.t. $|k| \leq q-1$. Using (14), we deduce that

$$
\begin{aligned}
B_{q} & =2 \sum_{r=0}^{q-1} m^{2(q-r-1)} \sum_{k \in \mathbb{G}_{r}} \mathbb{E}\left[P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)\left(X_{k}\right) \mathbf{1}_{\left\{k \in \mathbb{T}^{*}\right\}}\right] \\
& =2 \sum_{r=0}^{q-1} m^{2 q-r-2}\left\langle\nu, Q^{r} P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)\right\rangle
\end{aligned}
$$

Therefore, we get

$$
\begin{align*}
\left\|\tilde{M}_{\mathbb{G}_{q}^{*}}(f)\right\|_{L^{2}}^{2} & =m^{-2 q}\left\|\sum_{i \in \mathbb{G}_{q}^{*}} f\left(X_{i}\right)\right\|_{L^{2}}^{2} \\
& =m^{-q} \mathbb{E}\left[f^{2}\left(Y_{q}\right)\right]+2 m^{-2} \sum_{r=0}^{q-1} m^{-r}\left\langle v, Q^{r} P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)\right\rangle . \tag{15}
\end{align*}
$$

As $f \in F$, properties (ii), (iv), (v) and (vi) imply that $\lim _{q \rightarrow \infty} m^{-q} \mathbb{E}\left[f^{2}\left(Y_{q}\right)\right]=0$. Properties (iii), (iv) and (v) with $\langle\mu, f\rangle=0$ imply that $P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)$ converges to 0 as $q$ goes to infinity (with $r$ fixed) and is bounded uniformly in $q>r$ by a function of $F$. Thus, properties (v) and (vi) imply that $\left\langle v, Q^{r} P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)\right\rangle$ converges to 0 as $q$ goes to infinity (with $r$ fixed) and is bounded uniformly in $q>r$ by a finite constant, say $K$. For any $\varepsilon>0$, we can choose $r_{0}$ s.t. $\sum_{r>r_{0}} m^{-r} K \leq \varepsilon$ and $q_{0}>r_{0}$ s.t. for $q \geq q_{0}$ and $r \leq r_{0}$, we have $\left|\left\langle\nu, Q^{r} P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)\right\rangle\right| \leq \varepsilon / r_{0}$. We then get that for all $q \geq q_{0}$

$$
\sum_{r=0}^{q-1} m^{-r}\left|\left\langle\nu, Q^{r} P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)\right\rangle\right| \leq \sum_{r=0}^{r_{0}} r_{0}^{-1} \varepsilon+\sum_{r=r_{0}+1}^{q-1} m^{-r} K \leq 2 \varepsilon
$$

This gives that $\lim _{q \rightarrow \infty} \sum_{r=0}^{q-1} m^{-r}\left\langle v, Q^{r} P\left(Q^{q-r-1} f \otimes Q^{q-r-1} f\right)\right\rangle=0$. Finally, we get from (15) that if $\langle\mu, f\rangle=0$, then $\lim _{q \rightarrow \infty}\left\|\tilde{M}_{\mathbb{G}_{q}^{*}}(f)\right\|_{L^{2}}=0$.

For any function $f \in F$, we have, with $g=f-\langle\mu, f\rangle$,

$$
\tilde{M}_{\mathbb{G}_{q}^{*}}(f)=\tilde{M}_{\mathbb{G}_{q}^{*}}(g)+\langle\mu, f\rangle m^{-q}\left|\mathbb{G}_{q}^{*}\right| .
$$

As $g \in F$ and $\langle\mu, g\rangle=0$, the previous computations yield that $\lim _{q \rightarrow \infty}\left\|\tilde{M}_{\mathbb{G}_{q}^{*}}(g)\right\|_{L^{2}}=0$. As ( $m^{-q}\left|\mathbb{G}_{q}^{*}\right|, q \geq 1$ ) converges in $L^{2}$ (and a.s.) to $W$, we get that $\widetilde{M}_{\mathbb{G}_{q}^{*}}(f)$ converges to $\langle\mu, f\rangle W$ in $L^{2}$.

Then use that $m^{-q}\left|\mathbb{G}_{q}^{*}\right|$ converges a.s. to $W$ to get the second part of the Theorem.
We now prove a similar result for the average over the $r$ first generations. We set $t_{r}=\mathbb{E}\left[\left|\mathbb{T}_{r}^{*}\right|\right]$ and recall the explicit formulas given at (7). We first state an elementary Lemma, whose proof is left to the reader, and a second one which gives the asymptotic behavior of $\left|\mathbb{T}_{r}^{*}\right|$.

Lemma 3.3. Let ( $v_{r}, r \in \mathbb{N}$ ) be a sequence of real numbers converging to $a \in \mathbb{R}_{+}$, and $m$ a real such that $m>1$. Let

$$
w_{r}=\sum_{q=0}^{r} m^{q-r-1} v_{q} .
$$

Then the sequence $\left(w_{r}, r \in \mathbb{N}\right)$ converges to $a /(m-1)$.
The following Lemma is a direct consequence of Lemma 3.3 and of the definition of $W$ given by (8).

Lemma 3.4. We have $\lim _{q \rightarrow \infty} \frac{\left|\mathbb{G}_{q}^{*}\right|}{m^{q}}=\lim _{r \rightarrow \infty} \frac{\left|\mathbb{T}_{r}^{*}\right|}{t_{r}}=W$ a.s.
We now state the weak law of large numbers when averaging over all individuals up to a given generation.

Theorem 3.5. Let $\left(X_{i}, i \in \mathbb{T}^{*}\right)$ be a super-critical spatially homogeneous $P^{*}$-BMC on a $G W$ tree. Let $F$ satisfy (i)-(vi) and $f \in F$. The sequence $\left(\widetilde{M}_{\mathbb{T}_{r}^{*}}(f), r \in \mathbb{N}\right)$ converges to $\langle\mu, f\rangle W$ in $L^{2}$ as $r \rightarrow \infty$, where $W$ is defined by (8). We also have that the sequence $\left(\bar{M}_{\mathbb{T}_{r}^{*}}(f) \mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}}, r \in \mathbb{N}\right)$ converges to $\langle\mu, f\rangle \mathbf{1}_{\{W \neq 0\}}$ in probability when $r \rightarrow \infty$.

Proof. We have

$$
\begin{aligned}
\left\|\frac{1}{t_{r}} \sum_{i \in \mathbb{T}_{r}^{*}} f\left(X_{i}\right)-\langle\mu, f\rangle W\right\|_{L^{2}} & =\left\|\sum_{q=0}^{r} \frac{m^{q}}{t_{r}}\left(\tilde{M}_{\mathbb{G}_{q}^{*}}(f)-\langle\mu, f\rangle W\right)\right\|_{L^{2}} \\
& \leq \sum_{q=0}^{r} \frac{m^{q}}{t_{r}}\left\|\tilde{M}_{\mathbb{G}_{q}^{*}}(f)-\langle\mu, f\rangle W\right\|_{L^{2}} \\
& =\frac{m-1}{1-m^{-r-1}} \sum_{q=0}^{r} m^{q-r-1}\left\|\tilde{M}_{\mathbb{G}_{q}^{*}}(f)-\langle\mu, f\rangle W\right\|_{L^{2}} .
\end{aligned}
$$

The first part of the Theorem follows from Theorem 3.1 and Lemma 3.3. The second part of the Theorem is thus obtained using Lemma 3.4.

Remark 3.6. Applying Theorem 3.5 with $f=1$ ( $F$ contains the constants) immediately yields that $\left(t_{r}^{-1}\left|\mathbb{T}_{r}^{*}\right|, r \in \mathbb{N}\right)$ converges to $W$ in $L^{2}$ as $r$ goes to infinity.

The following Theorem extends those results to functions defined on the mother-daughters quantities of interest $\Delta_{i}=\left(X_{i}, X_{i 0}, X_{i 1}\right) \in \bar{S}^{3}$. Recall notations (6) and (9).

Theorem 3.7. Let $\left(X_{i}, i \in \mathbb{T}^{*}\right)$ be a super-critical spatially homogeneous $P^{*}$-BMC on a $G W$ tree. Let $F$ satisfy (i)-(vi) and $f \underset{\sim}{\mathcal{M}} \mathcal{B}\left(\bar{S}^{3}\right)$. We assume that $P^{*} f$ and $P^{*}\left(f^{2}\right)$ exist and belong to $F$. Then the sequences $\left(\widetilde{M}_{\mathbb{G}_{q}^{*}}(f), q \in \mathbb{N}\right)$ and $\left(\widetilde{M}_{\mathbb{T}_{r}^{*}}(f), r \in \mathbb{N}\right)$ converge to $\left\langle\mu, P^{*} f\right\rangle W$ in $L^{2}$, where $W$ is defined by (8); and the sequences $\left(\bar{M}_{\mathbb{G}_{q}^{*}}(f) \mathbf{1}_{\left\{\left|\mathbb{G}_{q}^{*}\right|>0\right\}}, q \in \mathbb{N}\right)$ and $\left(\bar{M}_{\mathbb{T}_{r}^{*}}(f) \mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}}, r \in \mathbb{N}\right)$ converge to $\left\langle\mu, P^{*} f\right\rangle \mathbf{1}_{\{W \neq 0\}}$ in probability.

Proof. Recall that $M_{\mathbb{G}_{q}^{*}}(f)=\sum_{i \in \mathbb{G}_{q}^{*}} f\left(\Delta_{i}\right)$. Let us compute $\left\|M_{\mathbb{G}_{q}^{*}}(f)\right\|_{L^{2}}^{2}$ :

$$
\left\|M_{\mathbb{G}_{q}^{*}}(f)\right\|_{L^{2}}^{2}=\sum_{i \in \mathbb{G}_{q}} \mathbb{E}\left[f^{2}\left(\Delta_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right]+\sum_{(i, j) \in \mathbb{G}_{q}^{2}, i \neq j} \mathbb{E}\left[f\left(\Delta_{i}\right) f\left(\Delta_{j}\right) \mathbf{1}_{\left\{(i, j) \in \mathbb{T}^{* 2}\right\}}\right]
$$

Remark that $\left\{i \in \mathbb{T}^{*}\right\}=\left\{X_{i} \neq \partial\right\}$, so that $\left\{i \in \mathbb{T}^{*}\right\}$ and $\left\{(i, j) \in \mathbb{T}^{* 2}\right\}$ both belong to $\sigma\left(X_{k}, k \in \mathbb{T}_{q}\right)$ for any $i, j$ in $\mathbb{G}_{q}$. We thus apply the Markov property for BMC to obtain

$$
\begin{aligned}
\left\|M_{\mathbb{G}_{q}^{*}}(f)\right\|_{L^{2}}^{2}= & \sum_{i \in \mathbb{G}_{q}} \mathbb{E}\left[P^{*}\left(f^{2}\right)\left(X_{i}\right) \mathbf{1}_{\left\{i \in \mathbb{T}^{*}\right\}}\right] \\
& +\sum_{(i, j) \in \mathbb{G}_{q}^{2}, i \neq j} \mathbb{E}\left[P^{*} f\left(X_{i}\right) P^{*} f\left(X_{j}\right) \mathbf{1}_{\left\{(i, j) \in \mathbb{T}^{*}{ }^{2}\right\}}\right] \\
= & \mathbb{E}\left[\left(\sum_{i \in \mathbb{G}_{q}^{*}} P^{*} f\left(X_{i}\right)\right)^{2}\right]-\mathbb{E}\left[\sum_{i \in \mathbb{G}_{q}^{*}}\left(P^{*} f\right)^{2}\left(X_{i}\right)\right] \\
& +\mathbb{E}\left[\sum_{i \in \mathbb{G}_{q}^{*}} P^{*}\left(f^{2}\right)\left(X_{i}\right)\right] \\
= & \left\|M_{\mathbb{G}_{q}^{*}}\left(P^{*} f\right)\right\|_{L^{2}}^{2}+\mathbb{E}\left[M_{\mathbb{G}_{q}^{*}}\left(P^{*}\left(f^{2}\right)-\left(P^{*} f\right)^{2}\right)\right] .
\end{aligned}
$$

Since $\left(m^{-q} M_{\mathbb{G}_{q}^{*}}\left(P^{*}\left(f^{2}\right)-\left(P^{*} f\right)^{2}\right), q \in \mathbb{N}\right)$ converges to $\left\langle\mu, P^{*}\left(f^{2}\right)-\left(P^{*} f\right)^{2}\right\rangle$ in $L^{2}$ and thus in $L^{1}$, we have that $m^{-2 q} \mathbb{E}\left[M_{\mathbb{G}_{q}^{*}}\left(P^{*}\left(f^{2}\right)-\left(P^{*} f\right)^{2}\right)\right]$ converges to 0 as $q$ goes to infinity. Then, we deduce the convergence of $\left(\tilde{M}_{\mathbb{G}_{q}^{*}}(f), q \in \mathbb{N}\right)$ and $\left(\bar{M}_{\mathbb{G}_{q}^{*}}(f) \mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}}, q \in \mathbb{N}\right)$ from Theorem 3.1.

The proof for the convergence of $\left(\tilde{M}_{\mathbb{T}_{r}^{*}}(f), r \in \mathbb{N}\right)$ and $\left(\bar{M}_{\mathbb{T}_{r}^{*}}(f) \mathbf{1}_{\left\{\left|\left|\mathbb{G}_{r}^{*}\right|>0\right\}\right.}, r \in \mathbb{N}\right)$ mimics then the proof of Theorem 3.5.

To end this section, we state strong laws of large numbers, under stronger assumptions on the auxiliary Markov chain $Y$.

Theorem 3.8. Let $\left(X_{i}, i \in \mathbb{T}^{*}\right)$ be a super-critical spatially homogeneous $P^{*}$-BMC on a $G W$ tree. Let $F$ satisfy (i)-(vi) and $f \in \mathcal{B}\left(\bar{S}^{3}\right)$. We assume that $P^{*} f$ and $P^{*}\left(f^{2}\right)$ exist and belong to $F$. We also suppose that there exist $c \in F$ and a non-negative sequence $\left(a_{r}, r \in \mathbb{N}\right.$ ) s.t. $\sum_{r \in \mathbb{N}} a_{r}^{2}<\infty$, and for all $x \in S$ and $r \in \mathbb{N},\left|Q^{r}\left(P^{*} f\right)(x)-\left\langle\mu, P^{*} f\right\rangle\right| \leq a_{r} c(x)$. Then the sequences $\left(\tilde{M}_{\mathbb{G}_{q}^{*}}(f), q \in \mathbb{N}\right)$ and $\left(\widetilde{M}_{\mathbb{T}_{r}^{*}}(f), r \in \mathbb{N}\right)$ converge to $\left\langle\mu, P^{*} f\right\rangle \bar{W}$ a.s., where $W$ is defined by (8); and the sequences $\left(\bar{M}_{\mathbb{G}_{q}^{*}}(f) \mathbf{1}_{\left\{\left|\mathbb{G}_{q}^{*}\right|>0\right\}}, q \in \mathbb{N}\right)$ and $\left(\bar{M}_{\mathbb{T}_{r}^{*}}(f) \mathbf{1}_{\left\{\left|\mathbb{G}_{r}^{*}\right|>0\right\}}, r \in \mathbb{N}\right)$ converge a.s. to $\left\langle\mu, P^{*} f\right\rangle \mathbf{1}_{\{W \neq 0\}}$.

Proof. Thanks to Lemma 3.4 and the equality

$$
\tilde{M}_{\mathbb{G}_{q}^{*}}(f)-\left\langle\mu, P^{*} f\right\rangle W=\tilde{M}_{\mathbb{G}_{q}^{*}}\left(f-\left\langle\mu, P^{*} f\right\rangle\right)+\left(m^{-q}\left|\mathbb{G}_{q}^{*}\right|-W\right)\left\langle\mu, P^{*} f\right\rangle,
$$

proving the a.s. convergence of $\left(\tilde{M}_{\mathbb{G}_{q}^{*}}(f), q \in \mathbb{N}\right)$ amounts to prove the a.s. convergence of $\left(\tilde{M}_{\mathbb{G}_{q}^{*}}(g), q \in \mathbb{N}\right)$ to 0 , where $g=f-\left\langle\mu, P^{*} f\right\rangle$. It is enough to show that $\sum_{q \geq 0}$ $\mathbb{E}\left[\left(\tilde{M}_{\mathbb{G}_{q}^{*}}(g)\right)^{2}\right]<\infty$. But we established in the proof of Theorem 3.7 that

$$
\mathbb{E}\left[\left(\tilde{M}_{\mathbb{G}_{q}^{*}}(g)\right)^{2}\right]=\left\|\tilde{M}_{\mathbb{G}_{q}^{*}}\left(P^{*} g\right)\right\|_{L^{2}}^{2}+m^{-2 q} \mathbb{E}\left[M_{\mathbb{G}_{q}^{*}}\left(P^{*}\left(g^{2}\right)-\left(P^{*} g\right)^{2}\right)\right]
$$

and since $P^{*} f$ and $P^{*}\left(f^{2}\right)$ both belong to $F$, we get the same for $P^{*} g$ and $P^{*}\left(g^{2}\right)$. We thus know that $m^{-q} \mathbb{E}\left[M_{\mathbb{G}_{q}^{*}}\left(P^{*}\left(g^{2}\right)-\left(P^{*} g\right)^{2}\right)\right]$ converges, so that $\sum_{q \geq 0} m^{-2 q} \mathbb{E}\left[M_{\mathbb{G}_{q}^{*}}\left(P^{*}\left(g^{2}\right)-\right.\right.$ $\left.\left.\left(P^{*} g\right)^{2}\right)\right]<\infty$. Finally, to obtain that $\sum_{q \geq 0}\left\|\tilde{M}_{\mathbb{G}_{q}^{*}}\left(P^{*} g\right)\right\|_{L^{2}}^{2}<\infty$, we follow the proofs of

Corollary 15 and Theorem 14 of [11], and the result is straightforward. Notice that the condition $\sum_{r \in \mathbb{N}} a_{r}<\infty$ of [11] can be weakened into $\sum_{r \in \mathbb{N}} a_{r}^{2}<\infty$.

Next, the a.s. convergence of ( $\widetilde{M}_{\mathbb{T}_{r}^{*}}(f), r \in \mathbb{N}$ ) is obtained from the previous one and Lemma 3.3. Finally, Lemma 3.4 allows to deduce the two last a.s. convergence from the previous ones.

## 4. Technical results about the weak law of large numbers

The technical Propositions of this Section deal with the average of a function $f$ when going through $\mathbb{T}^{*}$ via timescales $\left(\tau_{n}(t), t \in[0,1]\right)$ preserving the genealogical order. In order to define $\left(\tau_{n}(t), t \in[0,1]\right)$ we need to define $I_{n}^{*}$, set of the $n$ "first" cells of $\mathbb{T}^{*}$. Let ( $X_{i}, i \in \mathbb{T}^{*}$ ) be a super-critical spatially homogeneous $P^{*}$-BMC on a GW tree and $\mathcal{G}$ be the $\sigma$-field generated by ( $X_{i}, i \in \mathbb{T}$ ).

- We consider random variables $\left(\Pi_{q}^{*}, q \in \mathbb{N}^{*}\right)$ s.t. $\Pi_{q}^{*}$ takes values in the set of permutations of $\mathbb{G}_{q}^{*}$ for each $q$. Besides, conditionally on $\mathcal{G}$, these r.v. are independent and for all $q$, $\Pi_{q}^{*}$ is distributed as a uniform random permutation on $\mathbb{G}_{q}^{*}$. In particular, given $\left|\mathbb{G}_{q}^{*}\right|=k$, $\left(\Pi_{q}^{*}(1), \ldots, \Pi_{q}^{*}(k)\right)$ can be viewed as a random drawing of all the elements of $\mathbb{G}_{q}^{*}$, without replacement.
- For each integer $n \in \mathbb{N}^{*}$, we define the random variable $\rho_{n}=\inf \left\{k: n \leq\left|\mathbb{T}_{k}^{*}\right|\right\}$, with the convention $\inf \emptyset=\infty$. Loosely speaking, $\rho_{n}$ is the number of the generation to which belongs the $n$-th element of $\mathbb{T}^{*}$. Notice that $\rho_{1}=0$.
- Let $\tilde{\Pi}$ be the function from $\mathbb{N}^{*}$ to $\mathbb{T}^{*} \cup\left\{\partial_{\mathbb{T}}\right\}$, where $\partial_{\mathbb{T}}$ is a cemetery point added to $\mathbb{T}^{*}$, given by $\tilde{\Pi}(1)=\emptyset$ and for $k \geq 2$ :

$$
\tilde{\Pi}(k)= \begin{cases}\Pi_{\rho_{k}}^{*}\left(k-\left|\mathbb{T}_{\rho_{k}-1}^{*}\right|\right) & \text { if } \rho_{k}<+\infty \\ \partial_{\mathbb{T}} & \text { if } \rho_{k}=+\infty\end{cases}
$$

Notice that $\tilde{\Pi}$ defines a random order on $\mathbb{T}^{*}$ which preserves the genealogical order: if $k \leq n$ then $|\tilde{\Pi}(k)| \leq|\tilde{\Pi}(n)|$, with the convention $\left|\partial_{\mathbb{T}}\right|=\infty$. We thus define the set of the $n$ "first" elements of $\mathbb{T}^{*}$ (when $\left|\mathbb{T}^{*}\right| \geq n$ ):

$$
\begin{equation*}
I_{n}^{*}=\left\{\tilde{\Pi}(k), 1 \leq k \leq n \wedge\left|\mathbb{T}^{*}\right|\right\} \tag{16}
\end{equation*}
$$

We can now introduce the timescales: for $n \geq 1$, we consider the subdivision of $[0,1]$ given by $\left\{0, s_{n}, \ldots, s_{0}\right\}$, with $s_{k}=m^{-k}$. We define the continuous random time change $\left(\tau_{n}(t), t \in[0,1]\right)$ by

$$
\tau_{n}(t)= \begin{cases}m^{n} t, & t \in\left[0, m^{-n}\right],  \tag{17}\\ \left|\mathbb{T}_{n-k}^{*}\right|+\left(m^{k} t-1\right)(m-1)^{-1}\left|\mathbb{G}_{n-k+1}^{*}\right|, & t \in\left[m^{-k}, m^{-k+1}\right], 1 \leq k \leq n .\end{cases}
$$

Notice that $\tau_{n}(t) \leq\left|\mathbb{T}^{*}\right|$. The set $I_{\left\lfloor\tau_{n}(t)\right\rfloor}^{*}$, with $t \in[0,1]$, corresponds to the elements of $\mathbb{T}_{n-k}^{*}$, with $k=\left\lfloor-\frac{\log (t)}{\log (m)}\right\rfloor+1$, and the "first" fraction $\left(m^{k} t-1\right) /(m-1)$ of the elements of generation $\mathbb{G}_{n-k+1}^{*}$.

For the sake of simplicity, for any real $x \geq 0$, we will write $M_{x}^{*}(f)$ instead of $M_{I_{[x]}^{*}}(f)$ (recall (9)), with the convention that $M_{0}^{*}(f)=0$. We thus have, for $f \in \mathcal{B}\left(\bar{S}^{3}\right)$ e.g.,

$$
\begin{equation*}
M_{x}^{*}(f)=\sum_{i \in I_{\lfloor x\rfloor}^{*}} f\left(\Delta_{i}\right)=\sum_{k=1}^{\lfloor x\rfloor \wedge\left|\mathbb{T}^{*}\right|} f\left(\Delta_{\tilde{\Pi}(k)}\right) . \tag{18}
\end{equation*}
$$

Proposition 4.1. Let $F$ satisfy (i)-(vi) , $f \in F$ and $t \in[0,1]$. The sequence $\left(m^{-n} M_{\tau_{n}(t)}^{*}(f)\right.$, $n \in \mathbb{N}^{*}$ ) converges to $\langle\mu, f\rangle m(m-1)^{-1} W t$ in $L^{2}$ as $n$ goes to infinity.

Proof. We first consider the case $\langle\mu, f\rangle=0$. If $t=0$, then $\tau_{n}(t)=0$ and $M_{0}^{*}(f)=0$ by convention. Let $t \in(0,1]$ be fixed. There exists a unique $k \geq 1$ such that $m^{-k}<t \leq m^{-k+1}$. For $n \geq k$, we have, using (17) and that $\tilde{\Pi}$ preserves the order on $\mathbb{T}^{*}$,

$$
M_{\tau_{n}(t)}^{*}(f)=\sum_{i \in I_{\left\lfloor\tau_{n}(t)\right\rfloor}^{*}} f\left(X_{i}\right)=\sum_{i=1}^{\left\lfloor\tau_{n}(t)\right\rfloor} f\left(X_{\tilde{\Pi}(i)}\right)=M_{\mathbb{T}_{n-k}^{*}}(f)+M_{J_{n}}(f),
$$

where $J_{n}=\left\{\tilde{\Pi}(i),\left|\mathbb{T}_{n-k}^{*}\right|<i \leq\left\lfloor\tau_{n}(t)\right\rfloor\right\}$. Notice that $J_{n}=\emptyset$ if $\left|\mathbb{G}_{n-k+1}^{*}\right|=0$ and that, by convention, we then have $M_{J_{n}}(f)=0$. Both $k$ and $J_{n}$ depend on $t$, but since $t \in(0,1]$ is fixed, we shall not indicate this dependence. Theorem 3.5 implies that $m^{-n} M_{\mathbb{T}_{n-k}^{*}}(f)$ converges to 0 in $L^{2}$ as $n$ goes to $\infty$. Recall $\mathcal{G}$ is the $\sigma$-field generated by $\left(X_{i}, i \in \mathbb{T}\right)$. Since $J_{n} \subset \mathbb{G}_{n-k+1}^{*}$, we have

$$
\mathbb{E}\left[\left(M_{J_{n}}(f)\right)^{2} \mid \mathcal{G}\right]=\sum_{i, j \in \mathbb{G}_{n-k+1}^{*}} f\left(X_{i}\right) f\left(X_{j}\right) \mathbb{E}\left[\mathbf{1}_{\left\{i, j \in J_{n}\right\}} \mid \mathbb{T}^{*}\right] .
$$

Thanks to the definition of $\tilde{\Pi}$, we have for $i, j \in \mathbb{G}_{n-k+1}$

$$
\left.\mathbf{1}_{\left\{i, j \in \mathbb{G}_{n-k+1}^{*}\right\}} \mathbb{E}\left[\mathbf{1}_{\left\{i, j \in J_{n}\right\}}\right\} \mathbb{T}^{*}\right]=\mathbf{1}_{\left\{i, j \in \mathbb{G}_{n-k+1}^{*}\right\}}\left(\mathbf{1}_{\{i \neq j\}} \chi_{2}+\mathbf{1}_{\{i=j\}} \chi_{1}\right),
$$

where, with $a=\left\lfloor\left(m^{k} t-1\right)(m-1)^{-1}\left|\mathbb{G}_{n-k+1}^{*}\right|\right\rfloor$,

$$
\chi_{1}=\frac{a}{\left|\mathbb{G}_{n-k+1}^{*}\right|} \quad \text { and } \quad \chi_{2}=\frac{a(a-1)}{\left|\mathbb{G}_{n-k+1}^{*}\right|\left(\left|\mathbb{G}_{n-k+1}^{*}\right|-1\right)} .
$$

Thus, we get

$$
\begin{aligned}
\mathbb{E}\left[\left(M_{J_{n}}(f)\right)^{2} \mid \mathcal{G}\right] & =\chi_{2} \sum_{i, j \in \mathbb{G}_{n-k+1}^{*}} f\left(X_{i}\right) f\left(X_{j}\right)+\left(\chi_{1}-\chi_{2}\right) \sum_{i \in \mathbb{G}_{n-k+1}^{*}} f^{2}\left(X_{i}\right) \\
& =\chi_{2}\left(M_{\mathbb{G}_{n-k+1}^{*}}(f)\right)^{2}+\left(\chi_{1}-\chi_{2}\right) M_{\mathbb{G}_{n-k+1}^{*}}\left(f^{2}\right) \\
& \leq\left(M_{\mathbb{G}_{n-k+1}^{*}}(f)\right)^{2}+M_{\mathbb{G}_{n-k+1}^{*}}\left(f^{2}\right),
\end{aligned}
$$

as $0 \leq \chi_{2} \leq \chi_{1} \leq 1$. We thus have

$$
\begin{equation*}
\left\|m^{-n} M_{J_{n}}(f)\right\|_{L^{2}}^{2} \leq\left\|m^{-n} M_{\mathbb{G}_{n-k+1}^{*}}(f)\right\|_{L^{2}}^{2}+m^{-n}\left\|m^{-n} M_{\mathbb{G}_{n-k+1}^{*}}\left(f^{2}\right)\right\|_{L^{1}} \tag{19}
\end{equation*}
$$

The first term of the right-hand side of (19) converges to $\langle\mu, f\rangle W=0$ as $n$ goes to infinity, thanks to Theorem 3.1. The same Theorem entails that $\left\|m^{-n} M_{\mathbb{G}_{n-k+1}^{*}}\left(f^{2}\right)\right\|_{L^{1}}$ converges to $\mathbb{E}\left[\left\langle\mu, f^{2}\right\rangle W\right]$, and consequently the second term of the right-hand side of (19) also converges to 0 as $n$ goes to infinity. We deduce that the sequence ( $\left.m^{-n} M_{J_{n}}(f), m \in \mathbb{N}^{*}\right)$ converges to 0 in $L^{2}$.

Since $m^{-n} M_{\tau_{n}(t)}^{*}(f)=m^{-n} M_{\mathbb{T}_{n-k}^{*}}(f)+m^{-n} M_{J_{n}}(f)$, the sequence $\left(m^{-n} M_{\tau_{n}(t)}^{*}(f), n \in \mathbb{N}^{*}\right)$ converges to 0 in $L^{2}$.

Next, we consider the case $\langle\mu, f\rangle \neq 0$. We set $g=f-\langle\mu, f\rangle$. Since $m^{-n} M_{\tau_{n}(t)}^{*}(f)=$ $m^{-n} M_{\tau_{n}(t)}^{*}(g)+\langle\mu, f\rangle m^{-n}\left\lfloor\tau_{n}(t)\right\rfloor$, the Proposition will be proved as soon as we check that $\left(m^{-n}\left\lfloor\tau_{n}(t)\right\rfloor, n \in \mathbb{N}^{*}\right)$ converges to $m(m-1)^{-1} t W$ in $L^{2}$.

The case $t=0$ is obvious. For $t \in(0,1]$, there exists a unique $k \geq 1$ such that $m^{-k}<$ $t \leq m^{-k+1}$. We deduce from (17) that, for $1 \leq k \leq n$,

$$
m^{-n} \tau_{n}(t)=(m-1)^{-1}\left(\frac{\left|\mathbb{T}_{n-k}^{*}\right|}{t_{n-k}}\left(m^{-k+1}-\frac{1}{m^{n}}\right)+\frac{\left|\mathbb{G}_{n-k+1}^{*}\right|}{m^{n-k+1}}\left(m t-m^{-k+1}\right)\right)
$$

Since both $m^{-n}\left|\mathbb{G}_{n}^{*}\right|$ and $t_{n}^{-1}\left|\mathbb{T}_{n}^{*}\right|$ converges to $W$ in $L^{2}$ (see Remark 3.6), we finally obtain that $m^{-n} \tau_{n}(t)$ converges to $m(m-1)^{-1} t W$ in $L^{2}$.

We deduce from (7) and (17), that for $t \in(0,1], n \geq k$, where $k=\left\lfloor-\frac{\log (t)}{\log (m)}\right\rfloor+1$, we have

$$
\mathbb{E}\left[\tau_{n}(t)\right]=t_{n-k}+\left(m^{k} t-1\right)(m-1)^{-1} m^{n-k+1}=\left(m^{n+1} t-1\right)(m-1)^{-1} .
$$

Thus, Proposition 4.1 implies that $\left(\mathbb{E}\left[\tau_{n}(t)\right]^{-1} M_{\tau_{n}(t)}^{*}(f), n \in \mathbb{N}^{*}\right)$ converges to $\langle\mu, f\rangle W$ in $L^{2}$ for all $t \in[0,1]$, which generalizes Theorem 3.5.

In fact the convergence in Proposition 4.1 is uniform in $t$.
Corollary 4.2. Let $F$ satisfy (i)-(vi), $f \in F$ s.t. $|f| \in F$. We set $R_{n}(t)=m^{-n} M_{\tau_{n}(t)}^{*}(f)-$ $\langle\mu, f\rangle m(m-1)^{-1} W t$. The sequence $\left(\sup _{t \in[0,1]}\left|R_{n}(t)\right|, n \in \mathbb{N}^{*}\right)$ converges to 0 in $L^{2}$.

Proof. Let $f \in F$ s.t. $|f| \in F$. We set $f^{+}=\max (0, f)$ and $f^{-}=\max (0,-f)$. As $F$ is a vector space, we get that $f^{+}=(f+|f|) / 2$ and $f^{-}=f^{+}-f$ belong to $F$. Notice that $\left|R_{n}(t)\right| \leq\left|R_{n}^{+}(t)\right|+\left|R_{n}^{-}(t)\right|$, where $R_{n}^{\delta}(t)=m^{-n} M_{\tau_{n}(t)}^{*}\left(f^{\delta}\right)-\left\langle\mu, f^{\delta}\right\rangle m(m-1)^{-1} W t$ for $\delta \in\{+,-\}$. So it is enough to prove the Corollary for $f$ non-negative. As $t \mapsto m^{-n} M_{\tau_{n}(t)}^{*}(f)$ and $t \mapsto\langle\mu, f\rangle m(m-1)^{-1} t W$ are non-decreasing and $R_{n}(0)=0$, we get that for $N \geq 1$,

$$
\sup _{t \in[0,1]}\left|R_{n}(t)\right| \leq \frac{1}{N}\langle\mu, f\rangle m(m-1)^{-1} W+\sum_{k=1}^{N}\left|R_{n}(k / N)\right| .
$$

Now, use that $W \in L^{2}$ and that $R_{n}(t)$ goes to 0 in $L^{2}$ for all $t \in[0,1]$ to get the result.
We have a version of Proposition 4.1 and Corollary 4.2 for functions defined on $\bar{S}^{3}$.
Proposition 4.3. Let $F$ satisfy (i)-(vi), $g \in \mathcal{B}\left(\bar{S}^{3}\right)$ s.t. $P^{*} g$ and $P^{*}\left(g^{2}\right)$ exist and belong to $F$. Let $t \in[0,1]$. The sequence $\left(m^{-n} M_{\tau_{n}(t)}^{*}(g), n \in \mathbb{N}^{*}\right)$ converges to $\left\langle\mu, P^{*} g\right\rangle m(m-1)^{-1} t W$ in $L^{2}$.

Furthermore, if $P^{*}|g|$ and $P^{*}(g|g|)$ also belong to $F$ then $\left(\sup _{t \in[0,1]}\left|R_{n}(t)\right|, n \in \mathbb{N}^{*}\right)$ converges to 0 in $L^{2}$, where $R_{n}(t)=m^{-n} M_{\tau_{n}(t)}^{*}(g)-\left\langle\mu, P^{*} g\right\rangle m(m-1)^{-1} W t$, for $t \in[0,1]$.

Proof. The proof of the first part is similar to the proof of Theorem 3.7. The proof of the second part is similar to the proof of Corollary 4.2.

## 5. Fluctuations

Recalling (18), we shall prove a central limit theorem for the sequence $\left(M_{n}^{*}(f), n \geq 1\right)$, based on martingale theorems.

We set $\mathcal{H}_{n}=\sigma\left(\left(\Delta_{\tilde{\Pi}(k)}, 1 \leq k \leq n \wedge\left|\mathbb{T}^{*}\right|\right),(\tilde{\Pi}(k), 1 \leq k \leq n+1)\right)$ for $n \geq 1$, $\mathcal{H}_{0}=\sigma\left(X_{\emptyset}\right)$ and $\mathcal{H}=\left(\mathcal{H}_{n}, n \in \mathbb{N}\right)$ for the corresponding filtration. With the convention that $X_{\partial_{\mathbb{T}}}=\partial$, we notice that $X_{\tilde{\Pi}(n+1)}$ is $\mathcal{H}_{n}$-measurable. Indeed, given $(\tilde{\Pi}(k), 1 \leq k \leq n+1)$,
if $\tilde{\Pi}(n+1) \neq \partial_{\mathbb{T}}$, we have $\tilde{\Pi}(n+1)=\tilde{\Pi}(j) i$ for some $j \in\{1, \ldots, n\}$ and $i \in\{0,1\}$, and as $\Delta_{\tilde{\Pi}(j)}=\left(X_{\tilde{\Pi}(j)}, X_{\tilde{\Pi}(j) 0}, X_{\tilde{\Pi}(j) 1}\right) \in \mathcal{H}_{n}$, we deduce that $X_{\tilde{\Pi}(n+1)}$ is $\mathcal{H}_{n}$-measurable. In particular, as $\left\{\left|\mathbb{T}^{*}\right| \geq n+1\right\} \in \mathcal{H}_{n}$, this implies that $\mathbf{1}_{\left\{\left|\mathbb{T}^{*}\right| \geq n+1\right\}} \mathbb{E}\left[f\left(\Delta_{\tilde{\Pi}(n+1)}\right) \mid \mathcal{H}_{n}\right]=$ $\mathbf{1}_{\left\{\left|\mathbb{T}^{*}\right| \geq n+1\right\}} P^{*} f\left(X_{\tilde{\Pi}(n+1)}\right)$, for any $f \in \mathcal{B}\left(\bar{S}^{3}\right)$ such that $P^{*} f$ is well defined. If in addition $P^{*} f=0$, then $\left(M_{n}^{*}(f), n \in \mathbb{N}\right)$ is an $\mathcal{H}$-martingale.

We shall first recall a slightly weaker version of Theorem 4.3 from [19] on martingale convergence. (Theorem 4.3 from [19] is stated for filtrations which may vary with $n$.)

For $u \in \mathbb{R}^{d}$, we denote by $u^{\prime}$ its transpose. Let $\mathcal{H}=\left(\mathcal{H}_{i}, i \in \mathbb{N}\right)$ be a filtration. If ( $D_{i}, i \in \mathbb{N}$ ) is a sequence of vector-valued random variables $\mathcal{H}$-adapted and such that $\mathbb{E}\left[D_{i+1} \mid \mathcal{H}_{i}\right]=0$ for all $i \in \mathbb{N}$, then $\left(D_{i}, i \in \mathbb{N}\right)$ is called an $\mathcal{H}$-martingale difference.

Theorem 5.1 (Theorem 4.3 from [19]). Let $\mathcal{H}=\left(\mathcal{H}_{i}, i \in \mathbb{N}\right)$ be a filtration. For all $n \in \mathbb{N}^{*}$, let $\left(D_{n, i}=\left(D_{n, i}^{(1)}, \ldots, D_{n, i}^{(d)}\right)^{\prime}, i \in \mathbb{N}\right)$ be a sequence of $\mathbb{R}^{d}$-valued random vectors and an $\mathcal{H}$-martingale difference. For each $n \in \mathbb{N}$, let $\left(\tau_{n}(t), t \in[0,1]\right)$ be a non-decreasing càdlàg function s.t. $\tau_{n}(t)$ is a $\mathcal{H}$-stopping time for all $t \in[0,1]$. Let $(\mathcal{T}(t), t \in[0,1])$ be a $\mathbb{R}^{d \times d}$-valued continuous, possibly random, function. We assume the following two conditions hold:
(1) Convergence of the timescales. For all $t \in[0,1]$, we have the following convergence in probability:

$$
\sum_{i=1}^{\tau_{n}(t)} \mathbb{E}\left[D_{n, i}\left(D_{n, i}\right)^{\prime} \mid \mathcal{H}_{i-1}\right] \underset{n \rightarrow \infty}{\mathbb{P}} \mathcal{T}(t)
$$

(2) Lindeberg condition. For all $\varepsilon>0,1 \leq \ell \leq d$, we have the following convergence in probability:

$$
\sum_{i=1}^{\tau_{n}(1)} \mathbb{E}\left[\left(D_{n, i}^{(\ell)}\right)^{2} \mathbf{1}_{\left\{\left|D_{n, i}^{(\ell)}\right|>\varepsilon\right\}} \mid \mathcal{H}_{i-1}\right] \underset{n \rightarrow \infty}{\mathbb{P}} 0 .
$$

Then $\left(\sum_{i=1}^{\left\lfloor\tau_{n}(\cdot)\right\rfloor} D_{n, i}, n \in \mathbb{N}^{*}\right)$ converges in distribution to $B_{\mathcal{T}}$ in the Skorohod space $\mathbb{D}([0,1])^{\text {d }}$ of $\mathbb{R}^{d}$-valued càdlàg functions defined on $[0,1]$, where, conditionally on $\mathcal{T},\left(B_{\mathcal{T}}(t), t \geq 0\right)$ is a Gaussian process with independent increments and $B_{\mathcal{T}}(t)$ has zero mean and variance $\mathcal{T}(t)$.

Furthermore the convergence is stable: if $\left(Y_{n}, n \in \mathbb{N}\right)$ converges in probability to $Y$, then $\left(\left(\sum_{i=1}^{\left\lfloor\tau_{n} \cdot()\right\rfloor} D_{n, i}, Y_{n}\right), n \in \mathbb{N}\right)$ converges in distribution to $\left(B_{\mathcal{T}}, Y\right)$, where $B_{\mathcal{T}}$ is conditionally on $(\mathcal{T}, Y)$ distributed as $B_{\mathcal{T}}$ conditionally on $\mathcal{T}$, and the distribution of $(\mathcal{T}, Y)$ is determined by the following convergence

$$
\left(\sum_{i=1}^{\tau_{n}(\cdot)} \mathbb{E}\left[D_{n, i}\left(D_{n, i}\right)^{\prime} \mid \mathcal{H}_{i-1}\right], Y_{n}\right) \xrightarrow[n \rightarrow \infty]{\mathbb{P}}(\mathcal{T}, Y)
$$

We are now able to state the key result about fluctuations. For the sake of simplicity, we will write $P^{*} h^{k}$ for $P^{*}\left(h^{k}\right)$, and if $h=\left(h_{1}, \ldots, h_{d}\right)^{\prime}$ is an $\mathbb{R}^{d}$-valued function, we will write $P^{*} h$ for $\left(P^{*} h_{1}, \ldots, P^{*} h_{d}\right)^{\prime}$ and $\langle\mu, h\rangle$ for $\left(\left\langle\mu, h_{1}\right\rangle, \ldots,\left\langle\mu, h_{d}\right\rangle\right)^{\prime}$.

Theorem 5.2. Let $\left(X_{i}, i \in \mathbb{T}^{*}\right)$ be a super-critical spatially homogeneous $P^{*}$-BMC on a $G W$ tree and $\tau_{n}$ be defined by (17). Let $F$ satisfy (i)-(vi). Let $d \geq 1, d^{\prime} \geq 1, f=\left(f_{1}, \ldots, f_{d}\right)^{\prime} \in$
$\mathcal{B}\left(\bar{S}^{3}\right)^{d}, g=\left(g_{1}, \ldots, g_{d^{\prime}}\right)^{\prime} \in \mathcal{B}\left(\bar{S}^{3}\right)^{d^{\prime}}$ such that $P^{*} f_{\ell}^{k}$ exist and belong to $F$ for all $1 \leq \ell \leq d$ and $1 \leq k \leq 4, P^{*} g_{\ell}, P^{*}\left|g_{\ell}\right|, P^{*} g_{\ell}^{2}$ and $P^{*} g_{\ell}\left|g_{\ell}\right|$ exist and belong to $F$ for all $1 \leq \ell \leq d^{\prime}$. Let $\Sigma$ be a square root of the symmetric non-negative matrix $m(m-1)^{-1}\left\langle\mu, P^{*}\left(f f^{\prime}\right)-\left(P^{*} f\right)\left(P^{*} f\right)^{\prime}\right\rangle$ and $\gamma=m(m-1)^{-1}\left\langle\mu, P^{*} g\right\rangle$.

Then, the sequence $\left(m^{-n / 2} M_{\tau_{n}(\cdot)}^{*}\left(f-P^{*} f\right), m^{-n} M_{\tau_{n}(\cdot)}^{*}(g)\right)$ converges in distribution in the Skorohod space $\mathbb{D}\left([0,1], \mathbb{R}^{d+d^{\prime}}\right)$ of $\mathbb{R}^{d+d^{\prime}}$-valued càdlàg functions defined on $[0,1]$, to the process $\left(\Sigma \sqrt{W} B, \gamma W h_{0}\right)$, where $B$ is a d-dimensional Brownian motion independent of $W$, defined by (8), and $h_{0}$ is the identity function $t \mapsto t$.

Proof. Notice that $\tau_{n}$ defined by (17) is a non-decreasing continuous function s.t. $\tau_{n}(t)$ is a $\mathcal{H}$-stopping time for all $t \in[0,1]$. We set for all $n, i \in \mathbb{N}^{*}$,

$$
D_{n, i}=m^{-n / 2}\left(f\left(\Delta_{\tilde{\Pi}(i)}\right)-P^{*} f\left(X_{\tilde{\Pi}(i)}\right)\right) \mathbf{1}_{\left\{i \leq\left|\mathbb{T}^{*}\right|\right\}}
$$

so that $\left(D_{n, i}, i \in \mathbb{N}\right)$ is an $\mathcal{H}$-martingale difference. The matrix $\left\langle\mu, P^{*}\left(f f^{\prime}\right)-\left(P^{*} f\right)\left(P^{*} f\right)^{\prime}\right\rangle$ is indeed symmetric and non-negative, so that $\Sigma$ is well defined.

Notice that

$$
\begin{aligned}
\mathbb{E}\left[D_{n, i}\left(D_{n, i}\right)^{\prime} \mid \mathcal{H}_{i-1}\right]= & m^{-n}\left(P^{*}\left(f f^{\prime}\right)\left(X_{\tilde{\Pi}(i)}\right)\right. \\
& \left.-\left(P^{*} f\right)\left(X_{\tilde{\Pi}(i)}\right)\left(P^{*} f\right)^{\prime}\left(X_{\tilde{\Pi}(i)}\right)\right) \mathbf{1}_{\left\{i \leq\left|\mathbb{T}^{*}\right|\right\}}
\end{aligned}
$$

The convergence of the timescales (condition (1) of Theorem 5.1) to $\mathcal{T}(t)=\Sigma^{2} W t$, is then a direct application of Proposition 4.1

For $1 \leq \ell \leq d$, we have

$$
\begin{aligned}
\mathbb{E}\left[\left(D_{n, i}^{(\ell)}\right)^{2} \mathbf{1}_{\left\{\left|D_{n, i}^{(\ell)}\right|>\varepsilon\right\}} \mid \mathcal{H}_{i-1}\right] & \leq \varepsilon^{-2} \mathbb{E}\left[\left(D_{n, i}^{(\ell)}\right)^{4} \mid \mathcal{H}_{i-1}\right] \\
& =\varepsilon^{-2} m^{-2 n} P^{*}\left(f_{\ell}-P^{*} f_{\ell}\right)^{4}\left(X_{\tilde{\Pi}(i)}\right) \mathbf{1}_{\left\{i \leq\left|\mathbb{T}^{*}\right|\right\}}
\end{aligned}
$$

The Lindeberg condition of Theorem 5.1 is then a direct application of Proposition 4.1.
Notice the second part of Proposition 4.3. implies the convergence of $Y_{n}=m^{-n} M_{\tau_{n}(\cdot)}^{*}(g)$ to $\gamma W h_{0}$ in probability in the Skorohod space. We then deduce the result from Theorem 5.1.

The following result is an immediate consequence of Theorem 5.2.
Corollary 5.3. Let $\left(X_{i}, i \in \mathbb{T}^{*}\right)$ be a super-critical spatially homogeneous $P^{*}$-BMC on a $G W$ tree. Let $F$ satisfy (i)-(vi) . Let $f \in \mathcal{B}\left(\bar{S}^{3}\right)$ such that $P^{*} f^{k}$ exist and belong to $F$ for all $1 \leq k \leq 4$. Let $\sigma^{2}=\left\langle\mu, P^{*} f^{2}-\left(P^{*} f\right)^{2}\right\rangle$.

Then we have the following convergence in distribution:

$$
\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}\left|\mathbb{T}_{n}^{*}\right|^{-1 / 2} \sum_{i \in \mathbb{T}_{n}^{*}}\left(f\left(\Delta_{i}\right)-P^{*} f\left(X_{i}\right)\right) \xrightarrow[n \rightarrow \infty]{(d)} \mathbf{1}_{\{W \neq 0\}} \sigma G,
$$

where $G$ is a Gaussian random variable with mean zero and variance 1 independent of $W$, which is defined by (8).

Proof. Notice that $\sum_{i \in \mathbb{T}_{n}^{*}}\left(f\left(\Delta_{i}\right)-P^{*} f\left(X_{i}\right)\right)=M_{\tau_{n}(1)}^{*}(f)-M_{\tau_{n}(1)}^{*}\left(P^{*} f\right),\left|\mathbb{T}_{n}^{*}\right|=M_{\tau_{n}(1)}^{*}(\mathbf{1})$ and that $\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}$ converges a.s. to $\mathbf{1}_{\{W \neq 0\}}$. Then, to conclude, use the stable convergence of Theorem 5.2, and the fact that the marginals at time 1 converge since the limit is continuous.

The next result gives that fluctuations over different generations are asymptotically independent.

Corollary 5.4. Let $\left(X_{i}, i \in \mathbb{T}^{*}\right)$ be a super-critical spatially homogeneous $P^{*}$-BMC on a $G W$ tree. Let $F$ satisfy (i)-(vi). Let $d \geq 1, f_{1}, \ldots, f_{d} \in \mathcal{B}\left(\bar{S}^{3}\right)$ such that $P^{*} f_{\ell}^{k}$ exist and belong to $F$ for all $1 \leq \ell \leq d$ and $1 \leq k \leq 4$. Let $\sigma_{\ell}^{2}=\left\langle\mu, P^{*} f_{\ell}^{2}-\left(P^{*} f_{\ell}\right)^{2}\right\rangle$ for $1 \leq \ell \leq d$.

We set for $f \in \mathcal{B}\left(\bar{S}^{3}\right)$

$$
N_{n}(f)=\mathbf{1}_{\left\{\left|\left|\mathbb{G}_{n}^{*}\right|>0\right\}\right.}\left|\mathbb{G}_{n}^{*}\right|^{-1 / 2}\left(M_{\mathbb{G}_{n}^{*}}\left(f-P^{*} f\right)\right) .
$$

Then we have the following convergence in distribution:

$$
\left(N_{n}\left(f_{1}\right), \ldots, N_{n-d+1}\left(f_{d}\right)\right) \xrightarrow[n \rightarrow \infty]{(d)} \mathbf{1}_{\{W \neq 0\}}\left(\sigma_{1} G_{1}, \ldots, \sigma_{d} G_{d}\right),
$$

where $G_{1}, \ldots, G_{d}$ are independent Gaussian random variables with mean zero and variance 1 , and are independent of $W$, which is defined by (8).

Proof. Notice that for $n>k \geq 0$,

$$
N_{n-k}(f)=\frac{M_{\tau_{n}\left(m^{-k}\right)}^{*}\left(f-P^{*} f\right)-M_{\tau_{n}\left(m^{-k-1}\right)}^{*}\left(f-P^{*} f\right)}{\sqrt{M_{\tau_{n}\left(m^{-k}\right)}^{*}(\mathbf{1})-M_{\tau_{n}\left(m^{-k-1}\right)}^{*}(\mathbf{1})}} \mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}
$$

and $\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}$ converges a.s. to $\mathbf{1}_{\{W \neq 0\}}$. To conclude, use the stable convergence of Theorem 5.2 and that the increments of the Brownian motion are independent.

The extension of the two previous Corollaries to vector-valued functions can be proved in a very similar way. Notice also that the main results announced in Section 1.4 are by now established, in the light of Remark 2.2.

## 6. Estimation and tests for the asymmetric auto-regressive model

We consider the asymmetric auto-regressive model given in Section 1.1. Notice that the process $\left(X_{i}, i \in \mathbb{T}\right)$ defined in Section 1.1, with the convention that $X_{i}=\partial$ if the cell $i$ is dead or non existing, is a spatially homogeneous BMC on a GW tree. We shall assume it is super-critical, that is $2 p_{1,0}+p_{1}+p_{0}>1$.

We compute the maximum likelihood estimator (MLE)

$$
\hat{\theta}^{n}=\left(\hat{\alpha}_{0}^{n}, \hat{\beta}_{0}^{n}, \hat{\alpha}_{1}^{n}, \hat{\beta}_{1}^{n}, \hat{\alpha}_{0}^{\prime n}, \hat{\beta}_{0}^{\prime n}, \hat{\alpha}_{1}^{\prime n}, \hat{\beta}_{1}^{\prime n}, \hat{p}_{1,0}^{n}, \hat{p}_{0}^{n}, \hat{p}_{1}^{n}\right)
$$

of $\theta$ given by (4), and $\kappa^{n}=\left(\hat{\sigma}^{n}, \hat{\rho}^{n}, \hat{\sigma}_{0}^{n}, \hat{\sigma}_{1}^{n}\right)$ of $\kappa=\left(\sigma, \rho, \sigma_{0}, \sigma_{1}\right)$, based on the observation of a sub-tree $\mathbb{T}_{n+1}^{*}$. Let $\mathbb{T}_{n}^{1,0}$ be the set of cells in $\mathbb{T}_{n}^{*}$ with two living daughters, $\mathbb{T}_{n}^{0}$ (resp. $\mathbb{T}_{n}^{1}$ ) be the set of cells of $\mathbb{T}_{n}^{*}$ with only the new (resp. old) pole daughter alive:

$$
\begin{aligned}
& \mathbb{T}_{n}^{1,0}=\left\{i \in \mathbb{T}_{n}^{*}: \Delta_{i} \in S^{3}\right\}, \quad \mathbb{T}_{n}^{0}=\left\{i \in \mathbb{T}_{n}^{*}: \Delta_{i} \in S^{2} \times\{\partial\}\right\} \quad \text { and } \\
& \mathbb{T}_{n}^{1}=\left\{i \in \mathbb{T}_{n}^{*}: \Delta_{i} \in S \times\{\partial\} \times S\right\}
\end{aligned}
$$

It is elementary to get that for $\delta \in\{0,1\}$,

$$
\begin{align*}
& \hat{\alpha}_{\delta}^{n}=\frac{\left|\mathbb{T}_{n}^{1,0}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i} X_{i \delta}-\left(\left|\mathbb{T}_{n}^{1,0}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i}\right)\left(\left|\mathbb{T}_{n}^{1,0}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i \delta}\right)}{\left|\mathbb{T}_{n}^{1,0}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i}^{2}-\left(\left|\mathbb{T}_{n}^{1,0}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i}\right)^{2}}  \tag{20}\\
& \hat{\beta}_{\delta}^{n}=\left|\mathbb{T}_{n}^{1,0}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i \delta}-\hat{\alpha}_{\delta}^{n}\left|\mathbb{T}_{n}^{1,0}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i},  \tag{21}\\
& \hat{\alpha}_{\delta}^{\prime n}=\frac{\left|\mathbb{T}_{n}^{\delta}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i} X_{i \delta}-\left(\left|\mathbb{T}_{n}^{\delta}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i}\right)\left(\left|\mathbb{T}_{n}^{\delta}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i \delta}\right)}{\left|\mathbb{T}_{n}^{\delta}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i}^{2}-\left(\left|\mathbb{T}_{n}^{\delta}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i}\right)^{2}} \\
& \hat{\beta}_{\delta}^{\prime n}=\left|\mathbb{T}_{n}^{\delta}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i \delta}-\hat{\alpha}_{\delta}^{\prime n}\left|\mathbb{T}_{n}^{\delta}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i}, \\
& \hat{p}_{1,0}^{n}=\frac{\left|\mathbb{T}_{n}^{1,0}\right|}{\left|\mathbb{T}_{n}^{*}\right|}, \tag{22}
\end{align*} \hat{p}_{\delta}^{n}=\frac{\left|\mathbb{T}_{n}^{\delta}\right|}{\left|\mathbb{T}_{n}^{*}\right|}, ~ \$
$$

and

$$
\begin{aligned}
& \left(\hat{\sigma}^{n}\right)^{2}=\frac{1}{2\left|\mathbb{T}_{n}^{1,0}\right|} \sum_{i \in \mathbb{T}_{n}^{1,0}}\left(\hat{\varepsilon}_{i 0}^{2}+\hat{\varepsilon}_{i 1}^{2}\right), \quad \hat{\rho}^{n}=\frac{1}{\left(\hat{\sigma}^{n}\right)^{2}\left|\mathbb{T}_{n}^{1,0}\right|} \sum_{i \in \mathbb{T}_{n}^{1,0}} \hat{\varepsilon}_{i 0} \hat{\varepsilon}_{i 1}, \quad \text { and } \\
& \left(\hat{\sigma}_{\delta}^{n}\right)^{2}=\frac{1}{\left|\mathbb{T}_{n}^{\delta}\right|} \sum_{i \in \mathbb{T}_{n}^{\delta}} \hat{\varepsilon}_{i \delta}^{\prime 2}
\end{aligned}
$$

The residues are

$$
\hat{\varepsilon}_{i \delta}=X_{i \delta}-\hat{\alpha}_{\delta}^{n} X_{i}-\hat{\beta}_{\delta}^{n} \quad \text { for } i \in \mathbb{T}_{n}^{1,0}, \quad \text { and } \quad \hat{\varepsilon}_{i \delta}^{\prime}=X_{i \delta}-\hat{\alpha}_{\delta}^{\prime n} X_{i}-\hat{\beta}_{\delta}^{\prime n} \quad \text { for } i \in \mathbb{T}_{n}^{\delta}
$$

Notice that those MLE are based on polynomial functions of the observations. In order to use the results of Sections 3 and 5, we first show that the set of continuous and polynomially growing functions satisfies properties (i) to (v) of Section 2. The set of continuous and polynomially growing functions $C_{\mathrm{pol}}(\mathbb{R})$ is the set of continuous real functions defined on $\mathbb{R}$, satisfying that there exist $m \geq 0$ and $c \geq 0$ s.t. for all $x \in \mathbb{R},|f(x)| \leq c\left(1+|x|^{m}\right)$. It is easy to check that $C_{\text {pol }}(\mathbb{R})$ satisfies conditions (i)-(iii). To check properties (iv) and (v), we notice that the auxiliary Markov chain $Y=\left(Y_{n}, n \in \mathbb{N}\right)$ can be written in the following way:

$$
Y_{n+1}=a_{n+1} Y_{n}+b_{n+1}
$$

with $b_{n}=b_{n}^{\prime}+s_{n} e_{n}$, where $\left(\left(a_{n}, b_{n}^{\prime}, s_{n}\right), n \geq 1\right)$ is a sequence of independent identically distributed random variables, whose common distribution is given by

$$
\begin{equation*}
\mathbb{P}\left(a_{1}=\alpha_{\delta}, b_{1}^{\prime}=\beta_{\delta}, s_{1}=\sigma\right)=\frac{p_{1,0}}{m} \quad \text { and } \quad \mathbb{P}\left(a_{1}=\alpha_{\delta}^{\prime}, b_{1}^{\prime}=\beta_{\delta}^{\prime}, s_{1}=\sigma_{\delta}\right)=\frac{p_{\delta}}{m}, \tag{23}
\end{equation*}
$$

for $\delta \in\{0,1\}$. The sequence $\left(e_{n}, n \geq 1\right)$ is a sequence of independent $\mathcal{N}(0,1)$ random variables, and is independent of $\left(\left(a_{n}, b_{n}^{\prime}, s_{n}\right), n \geq 1\right)$, and both sequences are independent of $Y_{0}$. Notice that $Y_{n}$ is distributed as $Z_{n}=a_{1} a_{2} \cdots a_{n-1} a_{n} Y_{0}+\sum_{k=1}^{n} a_{1} a_{2} \cdots a_{k-1} b_{k}$. Since $\left|a_{k}\right| \leq \max \left(\left|\alpha_{0}\right|,\left|\alpha_{1}\right|,\left|\alpha_{0}^{\prime}\right|,\left|\alpha_{1}^{\prime}\right|\right)<1$ for all $k \in \mathbb{N}^{*}$, we get that the sequence $\left(Z_{n}, n \in \mathbb{N}\right)$ converges a.s. to a limit $Z$. This implies that $Y$ converges in distribution to $Z$. Following the proof of Lemma 26 in [11], we get that $C_{\mathrm{pol}}(\mathbb{R})$ fulfills properties (iv) and (v), with $\mu$ the distribution of $Z$.

Proposition 6.1. Assume that the distribution of the ancestor $X_{\emptyset}$ has finite moments of all orders. Then $\left(\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}} \hat{\theta}^{n}, n \geq 1\right)$ and $\left(\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}} \hat{\kappa}^{n}, n \geq 1\right)$ converge in probability respectively to $\mathbf{1}_{\{W \neq 0\}} \theta$ and $\mathbf{1}_{\{W \neq 0\}} \kappa$, where $W$ is defined by (8).

Proof. The hypothesis on the distribution of $X_{\emptyset}$ implies that $C_{\mathrm{pol}}(\mathbb{R})$ fulfills (vi). The result is then a direct consequence of Theorem 3.7.

Remark 6.2. Using similar arguments as in Propositions 30 and 34 of [11], it is easy to deduce from Theorem 3.8 and from the proof of Proposition 28 of [11], that the convergences in Proposition 6.1 hold a.s., that is the MLEs $\hat{\theta}^{n}$ and $\hat{\kappa}^{n}$ are strongly consistent.

From the definition of $Z_{n}$, we deduce that in distribution $Z \stackrel{(d)}{=} a_{1} Z^{\prime}+b_{1}$, where $Z^{\prime}$ is distributed as $Z$, and is independent of ( $a_{1}, b_{1}$ ) (see (23) for the distribution of $\left(a_{1}, b_{1}\right)$ ). This equality in distribution entails that

$$
\begin{equation*}
\mu_{1}=\mathbb{E}[Z]=\frac{\bar{\beta}}{1-\bar{\alpha}} \quad \text { and } \quad \mu_{2}=\mathbb{E}\left[Z^{2}\right]=\frac{2 \overline{\alpha \beta} \bar{\beta} /(1-\bar{\alpha})+\overline{\beta^{2}}+\overline{\sigma^{2}}}{1-\overline{\alpha^{2}}} \tag{24}
\end{equation*}
$$

where $\bar{\alpha}=\mathbb{E}\left[a_{1}\right], \overline{\alpha^{2}}=\mathbb{E}\left[a_{1}^{2}\right], \bar{\beta}=\mathbb{E}\left[b_{1}\right], \overline{\beta^{2}}=\mathbb{E}\left[b_{1}^{2}\right], \overline{\alpha \beta}=\mathbb{E}\left[a_{1} b_{1}\right]$ and $\overline{\sigma^{2}}=\mathbb{E}\left[s_{1}^{2}\right]$.
We can now state one of the main results of this section.
Proposition 6.3. Assume that the distribution of the ancestor $X_{\emptyset}$ has finite moments of all orders. Then $\left(\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}\left|\mathbb{T}_{n}^{*}\right|^{1 / 2}\left(\hat{\theta}^{n}-\theta\right), n \geq 1\right)$ converges in law to $\mathbf{1}_{\{W \neq 0\}} G_{11}$, where $G_{11}$ is a 11-dimensional vector, independent of $W$ defined by (8), of law $\mathcal{N}(0, \Sigma)$, with

$$
\begin{aligned}
\Sigma & =\left(\begin{array}{ccccc}
\sigma^{2} K / p_{1,0} & \rho \sigma^{2} K / p_{1,0} & 0 & 0 & 0 \\
\rho \sigma^{2} K / p_{1,0} & \sigma^{2} K / p_{1,0} & 0 & 0 & 0 \\
0 & 0 & \sigma_{0}^{2} K / p_{0} & 0 & 0 \\
0 & 0 & 0 & \sigma_{1}^{2} K / p_{1} & 0 \\
0 & 0 & 0 & 0 & \Gamma
\end{array}\right), \text { where } \\
K & =\left(\mu_{2}-\mu_{1}^{2}\right)^{-1}\left(\begin{array}{cc}
1 & -\mu_{1} \\
-\mu_{1} & \mu_{2}
\end{array}\right) \text { and } \\
\Gamma & =\left(\begin{array}{ccc}
p_{1,0}\left(1-p_{1,0}\right) & -p_{0} p_{1,0} & -p_{1} p_{1,0} \\
-p_{0} p_{1,0} & p_{0}\left(1-p_{0}\right) & -p_{0} p_{1} \\
-p_{1} p_{1,0} & -p_{0} p_{1} & p_{1}\left(1-p_{1}\right)
\end{array}\right)
\end{aligned}
$$

Remark 6.4. Notice that the eight first terms of $\theta$ are the parameters of the bifurcative autoregression model given by (1)-(3). In that framework, the matrix $\left(M(y)_{i, j}\right)_{(i, j) \in\{1, \ldots, 8\}^{2}}$, where $M(y)$ is defined in the proof below, is the inverse of $X^{\prime} X$ (up to a factor $\left|\mathbb{T}_{n}^{*}\right|$ ), where $X$ would be the design matrix of the bifurcative auto-regression, extending the notion defined in linear
regression. In the same way, the vector $U^{n}$ also defined in the proof below, coincides with $\left|\mathbb{T}_{n}^{*}\right|^{-1 / 2} X^{\prime} \varepsilon$, where $\varepsilon$ is the vector of error terms in the auto-regression. The formulas of standard linear regression giving the parameter's estimator and the difference between this estimator and the parameter, still hold in our framework of bifurcating auto-regression. Our proof thus consists in showing the convergence in probability of $\left|\mathbb{T}_{n}^{*}\right|\left(X^{\prime} X\right)^{-1}$, and that of $\left|\mathbb{T}_{n}^{*}\right|^{-1 / 2} X^{\prime} \varepsilon$ in distribution, with the stability of that last convergence.

Proof. This proof follows the idea of the proof of Proposition 33 of [11]. Let us introduce $\xi_{n}:=\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}\left|\mathbb{T}_{n}^{*}\right|^{1 / 2}\left(\hat{\theta}^{n}-\theta\right)$. We can rewrite $\xi_{n}$ as $\xi_{n}=\varphi\left(U^{n}, Y_{n}\right)$, where $\varphi(u, y)=M(y) u$ $\left(u \in \mathbb{R}^{11}, y \in \mathbb{R}^{9}\right.$ ), with $M(y)$ a matrix depending on $y$, defined by

$$
M(y)=\left(\begin{array}{ccccc}
N\left(y_{1}, y_{2}, y_{3}\right) & 0 & 0 & 0 & 0 \\
0 & N\left(y_{1}, y_{2}, y_{3}\right) & 0 & 0 & 0 \\
0 & 0 & N\left(y_{4}, y_{5}, y_{6}\right) & 0 & 0 \\
0 & 0 & 0 & N\left(y_{7}, y_{8}, y_{9}\right) & 0 \\
0 & 0 & 0 & 0 & I_{3}
\end{array}\right),
$$

where $I_{3}$ stands for the unit matrix of size 3 , and

$$
N(a, b, c)=\left(\begin{array}{cc}
c / b & -a / b \\
-a / b & \left(b+a^{2}\right) /(c b)
\end{array}\right) .
$$

As for the vectors $Y_{n}$ and $U^{n}$, the first one, $Y_{n}$, is a random vector of $\mathbb{R}^{9}$, defined by $Y_{n}:=\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}\left(A_{n}, B_{n}, C_{n}, A_{n}^{0}, B_{n}^{0}, C_{n}^{0}, A_{n}^{1}, B_{n}^{1}, C_{n}^{1}\right)$, where, for $\delta \in\{0,1\}$,

$$
\begin{aligned}
& A_{n}=\left|\mathbb{T}_{n}^{*}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i}, \\
& B_{n}=C_{n}\left(\left|\mathbb{T}_{n}^{*}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i}^{2}\right)-\left(\left|\mathbb{T}_{n}^{*}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{1,0}} X_{i}\right)^{2}, \\
& C_{n}=\left|\mathbb{T}_{n}^{*}\right|^{-1}\left|\mathbb{T}_{n}^{1,0}\right| \\
& A_{n}^{\delta}=\left|\mathbb{T}_{n}^{*}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i} \\
& B_{n}^{\delta}=C_{n}^{\delta}\left(\left|\mathbb{T}_{n}^{*}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i}^{2}\right)-\left(\left|\mathbb{T}_{n}^{*}\right|^{-1} \sum_{i \in \mathbb{T}_{n}^{\delta}} X_{i}\right)^{2} \\
& C_{n}^{\delta}=\left|\mathbb{T}_{n}^{*}\right|^{-1}\left|\mathbb{T}_{n}^{\delta}\right| .
\end{aligned}
$$

The second one, $U^{n}$, is a random vector of $\mathbb{R}^{11}$, defined by:

$$
U_{i}^{n}=\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}\left|\mathbb{T}_{n}^{*}\right|^{-1 / 2} \times \begin{cases}M_{\mathbb{T}_{n}^{*}}\left(f_{i}\right)-M_{\mathbb{T}_{n}^{1,0}}\left(P f_{i}\right) & \text { for } i \in\{1,2,3,4\}, \\ M_{\mathbb{T}_{n}^{*}}\left(f_{i}\right)-M_{\mathbb{T}_{n}^{0}}\left(P f_{i}\right) & \text { for } i \in\{5,6\}, \\ M_{\mathbb{T}_{n}^{*}}\left(f_{i}\right)-M_{\mathbb{T}_{n}^{1}}\left(P f_{i}\right) & \text { for } i \in\{7,8\}, \\ M_{\mathbb{T}_{n}^{*}}\left(f_{i}\right)-M_{\mathbb{T}_{n}^{*}}^{*}\left(P^{*} f_{i}\right) & \text { for } i \in\{9,10,11\},\end{cases}
$$

where $f_{1}=\mathbf{x y} 1_{S^{3}}, f_{2}=\mathbf{y 1}_{S^{3}}, f_{3}=\mathbf{x z} 1_{S^{3}}, f_{4}=\mathbf{z 1}_{S^{3}}, f_{5}=\mathbf{x y} 1_{S^{2} \times\{\partial\}}, f_{6}=\mathbf{y 1}_{S^{2} \times\{\partial\}}$, $f_{7}=\mathbf{x z} 1_{S \times\{\partial\} \times S}, f_{8}=\mathbf{z 1}_{S \times\{\partial\} \times S}, f_{9}=\mathbf{1}_{S^{3}}, f_{10}=\mathbf{1}_{S^{2} \times\{\partial\}}$ and $f_{11}=\mathbf{1}_{S \times\{\partial\} \times S}$.

The proof is thus organized as follows: in step 1, we prove the convergence in distribution of ( $U^{n}, n \geq 1$ ), and the stability of this convergence. In step 2 , we establish the convergence in probability of ( $Y_{n}, n \geq 1$ ). Finally, in step 3 , we deduce the result.

Step 1. In order to use Theorem 5.2, we write $U^{n}$ as $U^{n}=L V^{n}$, where $L$ is the matrix

$$
\begin{aligned}
L= & \left(\begin{array}{ccccc}
R\left(p_{1,0}\right) & 0 & 0 & 0 & 0 \\
0 & R\left(p_{1,0}\right) & 0 & 0 & 0 \\
0 & 0 & R\left(p_{0}\right) & 0 & 0 \\
0 & 0 & 0 & R\left(p_{1}\right) & 0 \\
0 & 0 & 0 & 0 & I_{3}
\end{array}\right), \\
& \text { with } R(p)=\left(\begin{array}{cccc}
1 & -1 / p & 0 & 0 \\
0 & 0 & 1 & -1 / p
\end{array}\right),
\end{aligned}
$$

and where $V^{n}$ is the random vector of $\mathbb{R}^{19}$, defined by

$$
V^{n}=\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}\left|\mathbb{T}_{n}^{*}\right|^{-1 / 2}\left(M_{\mathbb{T}_{n}^{*}}\left(g-P^{*} g\right)\right),
$$

with $g_{2 i-1}=f_{i}$ for $i \in\{1, \ldots, 8\}, g_{2 i}=\left(P^{*} f_{i}\right) \mathbf{1}_{S^{3}}$ for $i \in\{1, \ldots, 4\}, g_{2 i}=\left(P^{*} f_{i}\right) \mathbf{1}_{S^{2} \times\{\partial\}}$ for $i \in\{5,6\}, g_{2 i}=\left(P^{*} f_{i}\right) \mathbf{1}_{S \times\{\partial\} \times S}$ for $i \in\{7,8\}, g_{17}=\mathbf{1}_{S^{3}}, g_{18}=\mathbf{1}_{S^{2} \times\{\partial\}}$ and $g_{19}=\mathbf{1}_{S \times\{\partial\} \times S}$.
Now, Theorem 5.2 entails that $\left(m^{-n / 2} M_{\tau_{n}(\cdot)}^{*}\left(g-P^{*} g\right), n \geq 1\right)$ converges in distribution in the Skorohod space $\mathbb{D}\left([0,1], \mathbb{R}^{19}\right)$, to $\sqrt{m(m-1)^{-1} W} \Upsilon B$, where $\Upsilon$ is a square root of the symmetric non-negative matrix $\left\langle\mu, P^{*}\left(g g^{\prime}\right)-\left(P^{*} g\right)\left(P^{*} g\right)^{\prime}\right\rangle$, and $B$ a 19-dimensional Brownian motion independent of $W$. Besides, this convergence is stable. This immediately leads to the convergence of $\left(V^{n}, n \geq 1\right)$ to $\mathbf{1}_{\{W \neq 0\}} \Upsilon H$, where $H$ is a gaussian vector of law $\mathcal{N}\left(0, I_{19}\right)$. Indeed, the projection $f \mapsto f(1)$ is continuous on the Skorohod space, which gives the convergence in distribution of $m^{-n / 2} M_{\mathbb{T}_{n}^{*}}\left(g-P^{*} g\right)=m^{-n / 2} M_{\tau_{n}(1)}^{*}\left(g-P^{*} g\right)$. Next, since this convergence is still stable, and since $\left(\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}} m^{n / 2}\left|\mathbb{T}_{n}^{*}\right|^{-1 / 2}, n \geq 1\right)$ converges a.s. to $\mathbf{1}_{\{W \neq 0\}}\left(\sqrt{m(m-1)^{-1} W}\right)^{-1}$ (see Lemma 3.4), we get the convergence in distribution of ( $V^{n}, n \geq 1$ ). Notice that this convergence is still stable.

To close step 1 , recall that $U^{n}=L V^{n}$, so that the sequence ( $U^{n}, n \geq 1$ ) converges in law to $\mathbf{1}_{\{W \neq 0\}} G$, where $G$ is a centered gaussian vector of $\mathbb{R}^{19}$, independent of $W$, and with covariance matrix

$$
Q=\left(\begin{array}{ccccc}
\sigma^{2} p_{1,0} K^{-1} & \rho \sigma^{2} p_{1,0} K^{-1} & 0 & 0 & 0 \\
\rho \sigma^{2} p_{1,0} K^{-1} & \sigma^{2} p_{1,0} K^{-1} & 0 & 0 & 0 \\
0 & 0 & \sigma_{0}^{2} p_{0} K^{-1} & 0 & 0 \\
0 & 0 & 0 & \sigma_{1}^{2} p_{1} K^{-1} & 0 \\
0 & 0 & 0 & 0 & \Gamma
\end{array}\right) .
$$

This last convergence is once again stable.
Step 2. The convergence in probability of the sequence ( $Y_{n}, n \geq 1$ ) follows from Theorem 3.7, since each component of $Y_{n}$ is a continuous function of quantities $\bar{M}_{\mathbb{T}_{n}^{*}}(h)$, with $h$ functions of $\mathcal{B}\left(\bar{S}^{3}\right)$, such that $P^{*} h$ and $P^{*}\left(h^{2}\right)$ exist and belong to $F$, since $F$ is the set $C_{\mathrm{pol}}(\mathbb{R})$. Computations give that the limit of $\left(Y_{n}, n \geq 1\right)$ is $\mathbf{1}_{\{W \neq 0\}} \Lambda$, where $\Lambda$ is given by

$$
\Lambda=\left(p_{1,0} \mu_{1}, p_{1,0}^{2}\left(\mu_{2}-\mu_{1}^{2}\right), p_{1,0}, p_{0} \mu_{1}, p_{0}^{2}\left(\mu_{2}-\mu_{1}^{2}\right), p_{0}, p_{1} \mu_{1}, p_{1}^{2}\left(\mu_{2}-\mu_{1}^{2}\right), p_{1}\right)
$$

Step 3. Previous steps give that the sequence $\left(\left(U^{n}, Y_{n}\right), n \geq 1\right)$ converges in distribution to $\mathbf{1}_{\{W \neq 0\}}(G, \Lambda)$. Remind that $\xi_{n}=\varphi\left(U^{n}, Y_{n}\right)$, and since $\varphi$ is continuous, $\left(\xi_{n}, n \geq 1\right)$ converges in law to $\mathbf{1}_{\{W \neq 0\}} \varphi(G, \Lambda)$, so to speak to $\mathbf{1}_{\{W \neq 0\}} M(\Lambda) G$. Since $\Lambda$ is non-random, $G_{11}:=$ $M(\Lambda) G$ is a centered gaussian vector, independent of $W$, whose covariance matrix is given by $M(\Lambda) Q M(\Lambda)^{t}$. Notice that

$$
M(\Lambda)=\left(\begin{array}{ccccc}
K / p_{1,0} & 0 & 0 & 0 & 0 \\
0 & K / p_{1,0} & 0 & 0 & 0 \\
0 & 0 & K / p_{0} & 0 & 0 \\
0 & 0 & 0 & K / p_{1} & 0 \\
0 & 0 & 0 & 0 & I_{3}
\end{array}\right)
$$

which immediately gives that the covariance matrix of $G_{11}$ is $\Sigma$.
Remark 6.5. Proposition 6.3 deals with the asymptotic normality of the MLE of $\theta$ based on the observation of the sub-tree $\mathbb{T}_{n+1}^{*}$. If $L\left(X_{i}, i \in \mathbb{T}_{n+1}^{*}, \theta\right)$ denotes the corresponding log-likelihood function for $\theta$, the Fisher information, say $I_{n+1}$, is given by

$$
I_{n+1}=-\mathbb{E}\left[\frac{\partial^{2} L\left(X_{i}, i \in \mathbb{T}_{n+1}^{*}, \theta\right)}{\partial \theta \partial \theta^{\prime}}\right]
$$

Using Theorem 3.7, one can check that $\lim _{n \rightarrow \infty} I_{n+1} / \mathbb{E}\left[\left|\mathbb{T}_{n+1}^{*}\right|\right]=\Sigma^{-1}$. This is the analogue of the well-known asymptotic efficiency of the MLE for parametric sample of i.i.d. random variables.

Let $\theta_{1,0}\left(\right.$ resp. $\left.\hat{\theta}_{1,0}^{n}\right)$ stand for $\left(\alpha_{0}, \beta_{0}, \alpha_{1}, \beta_{1}\right)$ (resp. $\left.\left(\hat{\alpha}_{0}^{n}, \hat{\beta}_{0}^{n}, \hat{\alpha}_{1}^{n}, \hat{\beta}_{1}^{n}\right)\right)$.
Remark 6.6. Proposition 6.3 is quite similar to Proposition 33 in [11]. One of the main differences comes from the factor $p_{1,0}^{-1}$ in front of the matrix $K$ in the asymptotic covariance matrix for the estimation of $\theta_{1,0}$ with $\hat{\theta}_{1,0}^{n}$. As a matter of fact, this factor comes from the normalization by $\left|\mathbb{T}_{n}^{*}\right|^{1 / 2}$, number of living cells up to generation $n$, whereas this estimation is related to the cells with two living daughters, which would induce a normalization by $\left|\mathbb{T}_{n}^{1,0}\right|^{1 / 2}$. Since $\mathbf{1}_{\left\{\left|\mathbb{T}_{n}^{*}\right|>0\right\}}\left|\mathbb{T}_{n}^{1,0}\right| /\left|\mathbb{T}_{n}^{*}\right|$ converges in probability to $p_{1,0} \mathbf{1}_{\{W \neq 0\}}$, such a normalization would suppress the factor $p_{1,0}^{-1}$, see the following Corollary.

Corollary 6.7. Assume that the distribution of the ancestor $X_{\emptyset}$ has finite moments of all orders. Then $\left(\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}\left|\mathbb{T}_{n}^{1,0}\right|^{1 / 2}\left(\hat{\theta}_{1,0}^{n}-\theta_{1,0}\right), n \geq 1\right)$ converges in law to $\mathbf{1}_{\{W \neq 0\}} G_{4}$, where $G_{4}$ is a 4-dimensional vector, independent of $W$ defined by (8), with law $\mathcal{N}\left(0, \Sigma^{\prime}\right)$, with

$$
\Sigma^{\prime}=\sigma^{2}\left(\begin{array}{cc}
K & \rho K \\
\rho K & K
\end{array}\right), \quad \text { where } K=\left(\mu_{2}-\mu_{1}^{2}\right)^{-1}\left(\begin{array}{cc}
1 & -\mu_{1} \\
-\mu_{1} & \mu_{2}
\end{array}\right) .
$$

This way, this result is formally the same as Proposition 33 of [11], but one should notice that $\mu_{1}$ and $\mu_{2}$ are not defined the same way as in [11], since here they also depend on the parameters concerning cells with dead sisters. See Eqs. (2), (3), (23) and (24).

In order to detect cellular aging, see [11] in the case of no death ( $m=2$ ), we consider the null hypothesis $H_{0}=\left\{\left(\alpha_{0}, \beta_{0}\right)=\left(\alpha_{1}, \beta_{1}\right)\right\}$, which corresponds to no aging, and its alternative $H_{1}=\left\{\left(\alpha_{0}, \beta_{0}\right) \neq\left(\alpha_{1}, \beta_{1}\right)\right\}$. Notice that $\theta \mapsto \mu_{1}(\theta)$ and $(\theta, \kappa) \mapsto \mu_{2}(\theta, \kappa)$ given by (24) are continuous functions defined respectively on $\Theta=((-1,1) \times \mathbb{R})^{4} \times\left([0,1]^{3} \backslash\{0,0,0\}\right)$ and
$\Theta \times] 0,+\infty\left[{ }^{3}\right.$. We set $\hat{\mu}_{1}^{n}=\mu_{1}\left(\hat{\theta}^{n}\right)$ and $\hat{\mu}_{2}^{n}=\mu_{2}\left(\hat{\theta}^{n}, \hat{\kappa}^{n}\right)$. Proposition 6.8 allows to build a test for $H_{0}$ against $H_{1}$. Its proof, which is left to the reader, follows the proof of Proposition 35 of [11] and uses Corollary 6.7, the value of the extinction probability $\eta=\mathbb{P}(W=0)=1-\frac{m-1}{p_{1,0}}$, where $W$ is defined by (8), and Remark 6.2.

Proposition 6.8. Let $U$ and $V$ be two independent random variables, with $U$ distributed as a $\chi^{2}$ with two degrees of freedom, and V a Bernoulli random variable with parameter $1-\eta$.

Assume that the distribution of the ancestor $X_{\emptyset}$ has finite moments of all orders and define the test statistic

$$
\zeta_{n}=\frac{\left|\mathbb{T}_{n}^{1,0}\right|}{2\left(\hat{\sigma}^{n}\right)^{2}\left(1-\hat{\rho}^{n}\right)}\left\{\left(\hat{\alpha}_{0}^{n}-\hat{\alpha}_{1}^{n}\right)^{2}\left(\hat{\mu}_{2}^{n}-\left(\hat{\mu}_{1}^{n}\right)^{2}\right)+\left(\left(\hat{\alpha}_{0}^{n}-\hat{\alpha}_{1}^{n}\right) \hat{\mu}_{1}^{n}+\hat{\beta}_{0}^{n}-\hat{\beta}_{1}^{n}\right)^{2}\right\}
$$

Then, the statistics $\left(\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}} \zeta_{n}, n \geq 1\right)$ converges under $H_{0}$ in distribution to $U V$, and under $H_{1}$ a.s. to 0 on $\{V=0\}$ and $+\infty$ on $\{V=1\}$.

Remark 6.9. Let us assume that:

- Death occurs, that is $m \in(1,2)$.
- There is no difference for the marginal distribution of a daughter according to her sister is dead or alive; that is $\alpha_{\delta}^{\prime}=\alpha_{\delta}$ and $\beta_{\delta}^{\prime}=\beta_{\delta}$ for $\delta \in\{0,1\}$.
- For simplicity, the death probability is symmetric, that is $p_{0}=p_{1}$.

If one uses the statistics given by Proposition 33 in [11] with all the available data, that is if one uses

- Formula (20) and (21) with $\mathbb{T}_{n}^{1,0}$ replaced by $\mathbb{T}_{n}^{1,0} \cup \mathbb{T}_{n}^{\delta}$;
- The variance estimator:

$$
\left(\hat{\sigma}^{n}\right)^{2}=\frac{1}{\left|\mathbb{T}_{n+1}^{*}\right|-1}\left(\sum_{i \in \mathbb{T}_{n}^{1,0}}\left(\hat{\varepsilon}_{i 0}^{2}+\hat{\varepsilon}_{i 1}^{2}\right)+\sum_{i \in \mathbb{T}_{n}^{0}} \hat{\varepsilon}_{i 0}^{2}+\sum_{i \in \mathbb{T}_{n}^{1}} \hat{\varepsilon}_{i 1}^{2}\right)
$$

(Notice that we divide by $\left|\mathbb{T}_{n+1}^{*}\right|-1$ as this is equal to the total number of data: $2\left|\mathbb{T}_{n}^{1,0}\right|+$ $\left|\mathbb{T}_{n}^{0}\right|+\left|\mathbb{T}_{n}^{1}\right|$.)

- Keep the same estimation of the correlation: $\hat{\rho}^{n}=\frac{1}{\left(\hat{\sigma}^{n}\right)^{2}\left|\mathbb{T}_{n}^{1,0}\right|} \sum_{i \in \mathbb{T}_{n}^{1,0}} \hat{\varepsilon}_{i 0} \hat{\varepsilon}_{i 1}$;
then one checks that, as $n$ goes to infinity, $\mathbf{1}_{\left\{\left|\mathbb{G}_{n}^{*}\right|>0\right\}}\left|\mathbb{T}_{n}^{*}\right|^{1 / 2}\left(\hat{\theta}^{n}-\theta\right)$ converges in distribution to $\mathbf{1}_{\{W \neq 0\}} G$, where $G$ is a centered Gaussian vector with covariance matrix

$$
\sigma^{2}\left(p_{1,0}+p_{1}\right)^{-1}\left(\begin{array}{cc}
K & \rho p_{1,0}\left(p_{1,0}+p_{1}\right)^{-1} K \\
\rho p_{1,0}\left(p_{1,0}+p_{1}\right)^{-1} K & K
\end{array}\right)
$$

where $K$ is as in Proposition 6.3; and $G$ is independent of $W$, which is defined by (8). Then, it is not difficult to check that the statistics proposed by Guyon in Proposition 35 of [11], converges under $H_{0}$ towards $c U V$, with $U$ and $V$ as in Proposition 6.8 and

$$
c=\frac{\left(p_{1,0}+p_{1}\right)^{-1}\left(1-\rho p_{1,0}\left(p_{1,0}+p_{1}\right)^{-1}\right)}{(1-\rho)}
$$

As $\rho \in[-1,1], p_{1,0}+p_{1}>1 / 2$ (because $m>1$ and $p_{0}=p_{1}$ ), and $2 p_{1}+p_{1,0} \leq 1$, one can check that $c>1$. In particular, using the test statistic designed for cells with no death to data of cells with death leads to a non-conservative test.

## References

[1] M. Ackermann, L. CHAO, C.T. Bergstrom, M. Doebeli, On the evolutionary of aging, Aging Cell 6 (2) (2007) 235-244.
[2] M. Ackermann, S. Stearns, U. Jenal, Senescence in a bacterium with asymmetric division, Science 300 (2003) 1920.
[3] K.B. Athreya, H.-J. Kang, Some limit theorems for positive recurrent branching Markov chains, I, II, Adv. Appl. Probab. 30 (3) (1998) 693-710, 711-722.
[4] K.B. Athreya, P.E. Ney, Branching Processes, Springer-Verlag, New York, 1972.
[5] V. Bansaye, Proliferating parasites in dividing cells:Kimmel's branching model revisited, Ann. Appl. Probab. 18 (3) (2008) 967-996.
[6] V. Bansaye, Cell contamination and branching process in random environment with immigration, Adv. Appl. Probab. 41 (4) (2009) 1059-1081.
[7] V. Bansaye, J.-F. Delmas, L. Marsalle, V.C. Tran, Limit theorems for Markov processes indexed by continuous time Galton-Watson trees, Hal. (2009) 00431118v2.
[8] B. Bercu, B. De Saporta, A. Gégout-Petit, Asymptotic analysis for bifurcating autoregressive processes via a martingale approach, Electron J. Probab. 14 (87) (2009) 2492-2526.
[9] R. Cowan, R.G. Staudte, The bifurcating autoregression model in cell lineage, Biometrics 42 (4) (1986) 769-783.
[10] S.N. Evans, D. Steinsaltz, Damage segregation at fissioning may increase growth rates: a superprocess model, Theor. Pop. Biol. 71 (4) (2007) 473-490.
[11] J. Guyon, Limit theorems for bifurcating Markov chains. Application to the detection of cellular aging, Ann. Appl. Probab. 17 (5-6) (2007) 1538-1569.
[12] T. Harris, The Theory of Branching Process, Springer-Verlag, Berlin, 1963.
[13] H. Huang, W. Yang, Strong law of large numbers for Markov chains indexed by an infinite tree with uniformly bounded degree, Sci. China Ser. A 51 (2) (2008) 195-202.
[14] S.Y. Hwang, I.V. Basawa, I.K. Yeo, Local asymptotic normality for bifurcating autoregressive processes and related asymptotic inference, Stat. Methodol. 6 (2009) 61-69.
[15] A.B. Lindner, R. Madden, A. Demarez, E.J. Stewart, F. Taddei, Asymmetric segregation of protein aggregates is associated with cellular aging and rejuvenation, PNAS 105 (8) (2008) 3076-3081.
[16] F. Maaouia, A. Touati, Identification of multitype branching processes, C. R. Acad. Sci. Series I Math. 331 (2000) 923-928.
[17] S.P. Meyn, R.L. Tweedie, Markov chains and stochastic stability, in: Communications and Control Engineering Series, Springer-Verlag London Ltd., London, 1993.
[18] R. Pemantle, Tree-indexed processes, Statist. Sci. 10 (2) (1995) 200-213.
[19] H. Rootzén, Central limit theory for martingales via random change of time, in: Probability and Mathematical Statistics, Uppsala Univ., Uppsala, 1983. pp. 154-189.
[20] E.J. Stewart, R. Madden, G. Paul, F. Taddei, Aging and death in an organism that reproduces by morphologically symmetric division, PLoS Biol. 3 (2) (2005) e45.
[21] W. Yang, Some limit properties for Markov chains indexed by a homogeneous tree, Statist. Probab. Lett. 65 (3) (2003) 241-250.


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