# Prior sensitivity analysis in a semi-parametric <br> integer-valued time series model 

Helton Graziadei Antonio Lijoi<br>USP - São Paulo<br>Bocconi - Milano

| Hedibert F. Lopes | Paulo C. Marques F. | Igor Prünster |
| :---: | :---: | :---: |
| Insper - São Paulo | Insper - São Paulo | Bocconi - Milano |

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#### Abstract

We examine issues of prior sensitivity in a semi-parametric hierarchical extension of the $\operatorname{INAR}(p)$ model with innovation rates clustered according to a Pitman-Yor process placed at the top of the model hierarchy. Our main finding is a graphical criterion which guides the specification of the hyperparameters of the Pitman-Yor process base measure. It is shown how the discount and concentration parameters interact with the chosen base measure to yield a gain in terms of the robustness of the inferential results.


## 1 Introduction

Integer-valued time series are relevant to many fields of knowledge, ranging from finance and econometrics to ecology and meteorology. An extensive number of models for this kind of data have been proposed since the introduction of the $\operatorname{INAR}(1)$ model in the pioneering works of McKenzie [1] and Al-Osh and Alzaid [2] (see also the book by Weiss [3]). A higher-order $\operatorname{INAR}(p)$ model was considered in the work of Du and Li [4].

In this paper, we generalize the Bayesian version of the $\operatorname{INAR}(p)$ model studied by Neal and Kypraios [5]. In our model, the innovation rates are allowed to vary through time, with the distribution of the innovation rates being modeled hierarchically by means of a Pitman-Yor process [6]. In this way, we account for potential heterogeneities in the innovation rates as the process evolves through time. The semi-parametric form of the model demands a robustness analysis of our inferential conclusions as we vary the hyperparameters of the Pitman-Yor process. We investigate this prior sensitivity issue carefully and find ways to control the hyperparameters in order to achieve robust results.

The paper is organized as follows. In Section 2, we construct a generalized $\operatorname{INAR}(p)$ model with variable innovation rates. The likelihood function of the generalized model is derived and a data augmentation scheme is developed, which gives a specification of the model in terms of conditional distributions. This data augmented representation of the model enables the derivation in Section 4 of full conditional distributions in simple analytical form, which are essential for the stochastic simulations in Section 5. Section 3 recollects the main properties of the Pitman-Yor process which are used to define the PY-INAR $(p)$ model in Section 4, including its clustering properties. In building the $\operatorname{PY}-\operatorname{INAR}(p)$, we propose a form for the prior distribution of the thinning parameters vector which improves on the choice made for the Bayesian $\operatorname{INAR}(p)$ model studied in [5]. In Section 5, we investigate the robustness of the inference with respect to changes in the Pitman-Yor process hyperparameters. Using the full conditional distributions of the innovation rates derived in Section 4, we inspect the behavior of the model as we concentrate or spread the mass of the Pitman-Yor base measure. This leads us to a graphical criterion which identifies an elbow in the posterior expectation of the number of clusters as we vary the hyperparameters of the base measure. Once we have control over the base measure, we study its interaction with the concentration and discount hyperparameters, showing how to make choices that yield robust results. In the course of this development we use geometrical tools to inspect the clustering of the innovation rates produced by the model.

## 2 A generalization of the $\operatorname{INAR}(p)$ model

We begin by generalizing the original $\operatorname{INAR}(p)$ model of Du and $\mathrm{Li}[4]$ as follows.
Let $\left\{Y_{t}\right\}_{t \geq 1}$ be an integer-valued time series, and, for some integer $p \geq 1$, let the innovations $\left\{Z_{t}\right\}_{t \geq p+1}$, given positive parameters $\left\{\lambda_{t}\right\}_{t \geq p+1}$, be a sequence of conditionally independent Poisson $\left(\lambda_{t}\right)$ random variables. For a given vector of parameters $\alpha=$ $\left(\alpha_{1}, \ldots, \alpha_{p}\right) \in[0,1]^{p}$, let $\mathscr{F}_{i}=\left\{B_{i j}(t): j \geq 0, t \geq 2\right\}$ be a family of conditionally independent and identically distributed $\operatorname{Bernoulli}\left(\alpha_{i}\right)$ random variables. For $i \neq k$, suppose that $\mathscr{F}_{i}$ and $\mathscr{F}_{k}$ are conditionally independent, given $\alpha$. Furthermore, assume that the innovations $\left\{Z_{t}\right\}_{t \geq p+1}$ and the families $\mathscr{F}_{1}, \ldots, \mathscr{F}_{p}$ are conditionally independent, given $\alpha$ and $\lambda$. The generalized $\operatorname{INAR}(p)$ model is defined by the functional relation

$$
Y_{t}=\alpha_{1} \circ Y_{t-1}+\cdots+\alpha_{p} \circ Y_{t-p}+Z_{t}
$$

for $t \geq p+1$, in which $\circ$ denotes the binomial thinning operator, defined by $\alpha_{i} \circ Y_{t-i}=$ $\sum_{j=1}^{Y_{t-i}} B_{i j}(t)$, if $Y_{t-i}>0$, and $\alpha_{i} \circ Y_{t-i}=0$, if $Y_{t-i}=0$. In the homogeneous case, when all the $\lambda_{t}$ 's are assumed to be equal, we recover the original $\operatorname{INAR}(p)$ model.

When $p=1$, this model can be interpreted as specifying a birth-and-death process, in which, at epoch $t$, the number of cases $Y_{t}$ is equal to the new cases $Z_{t}$ plus the cases that survived from the previous epoch; the role of the binomial thinning operator being
to remove a random number of the $Y_{t-1}$ cases present at the previous epoch $t-1$ (see [7] for an interpretation of the order $p$ case as a birth-and-death process with immigration).

Let $y=\left(y_{1}, \ldots, y_{T}\right)$ denote the values of an observed time series. For simplicity, we assume that $Y_{1}=y_{1}, \ldots, Y_{p}=y_{p}$ with probability one. The joint distribution of $Y_{1}, \ldots, Y_{T}$, given parameters $\alpha$ and $\lambda=\left(\lambda_{p+1}, \ldots, \lambda_{T}\right)$, can be factored as

$$
\operatorname{Pr}\left\{Y_{1}=y_{1}, \ldots, Y_{T}=y_{T} \mid \alpha, \lambda\right\}=\prod_{t=p+1}^{T} \operatorname{Pr}\left\{Y_{t}=y_{t} \mid Y_{t-1}=y_{t-1}, \ldots, Y_{t-p}=y_{t-p}, \alpha, \lambda_{t}\right\}
$$

Since, with probability one, $\alpha_{i} \circ Y_{t-i} \leq Y_{t-i}$ and $Z_{t} \geq 0$, the likelihood function of the generalized $\operatorname{INAR}(p)$ model is given by

$$
\begin{aligned}
& L_{y}(\alpha, \lambda)=\prod_{t=p+1}^{T} \sum_{m_{1, t}=0}^{\min \left\{y_{t}, y_{t-1}\right\}} \cdots \sum_{m_{p, t}=0}^{\min \left\{y_{t}-\sum_{j=1}^{p-1} m_{j, t}, y_{t-p}\right\}}\left(\prod_{i=1}^{p}\binom{y_{t-i}}{m_{i, t}} \alpha_{i}^{m_{i, t}}\left(1-\alpha_{i}\right)^{y_{t-i}-m_{i, t}}\right) \times \\
&\left(\frac{e^{-\lambda_{t}} \lambda_{t}^{y_{t}-\sum_{j=1}^{p} m_{j, t}}}{\left(y_{t}-\sum_{j=1}^{p} m_{j, t}\right)!}\right)
\end{aligned}
$$

For some epoch $t$ and $i=1, \ldots, p$, suppose that we could observe the values of the latent maturations $M_{i, t}$. Postulate that

$$
M_{i, t} \mid Y_{t-i}=y_{t-i}, \alpha_{i} \sim \operatorname{Binomial}\left(y_{t-i}, \alpha_{i}\right)
$$

so that the conditional probability function of $M_{i, t}$ is given by

$$
\begin{aligned}
p\left(m_{i, t} \mid y_{t-i}, \alpha_{i}\right) & =\operatorname{Pr}\left\{M_{i, t}=m_{i, t} \mid Y_{t-i}=y_{t-i}, \alpha_{i}\right\} \\
& =\binom{y_{t-i}}{m_{i, t}} \alpha_{i}^{m_{i, t}}\left(1-\alpha_{i}\right)^{y_{t-i}-m_{i, t}} \mathbb{I}_{\left\{0, \ldots, y_{t-i}\right\}}\left(m_{i, t}\right)
\end{aligned}
$$

Furthermore, suppose that

$$
\begin{aligned}
p\left(y_{t} \mid m_{1, t}, \ldots, m_{p, t}, \lambda_{t}\right) & =\operatorname{Pr}\left\{Y_{t}=y_{t} \mid M_{1, t}=m_{1, t}, \ldots, M_{p, t}=m_{p, t}, \lambda_{t}\right\} \\
& =\frac{e^{-\lambda_{t}} \lambda_{t}^{y_{t}-\sum_{j=1}^{p} m_{j, t}}}{\left(y_{t}-\sum_{j=1}^{p} m_{j, t}\right)!} \mathbb{I}_{\left\{\sum_{j=1}^{p} m_{j, t}, \sum_{j=1}^{p} m_{j, t}+1, \ldots\right\}}\left(y_{t}\right)
\end{aligned}
$$

Using the law of total probability and the product rule, we have that

$$
\begin{aligned}
p\left(y_{t} \mid y_{t-1}, \ldots, y_{t-p}, \alpha, \lambda_{t}\right) & =\sum_{m_{1, t}=0}^{y_{t-1}} \cdots \sum_{m_{p, t}=0}^{y_{t-p}} p\left(y_{t}, m_{1, t}, \ldots, m_{p, t} \mid y_{t-1}, \ldots, y_{t-p}, \alpha, \lambda_{t}\right) \\
& =\sum_{m_{1, t}=0}^{y_{t-1}} \cdots \sum_{m_{p, t}=0}^{y_{t-p}} p\left(y_{t} \mid m_{1, t}, \ldots, m_{p, t}, \lambda_{t}\right) \times \prod_{i=1}^{p} p\left(m_{i, t} \mid y_{t-i}, \alpha_{i}\right)
\end{aligned}
$$

Since

$$
\begin{aligned}
\mathbb{I}_{\left\{\sum_{j=1}^{p} m_{j, t}, \sum_{j=1}^{p} m_{j, t}+1, \ldots\right\}}\left(y_{t}\right) & =\mathbb{I}_{\left\{0, \ldots, y_{t}\right\}}\left(\sum_{j=1}^{p} m_{j, t}\right) \\
& =\mathbb{I}_{\left\{0, \ldots, y_{t}\right\}}\left(m_{1, t}\right) \times \cdots \times \mathbb{I}_{\left\{0, \ldots, y_{t}-\sum_{j=1}^{p-1} m_{j, t}\right\}}\left(m_{p, t}\right)
\end{aligned}
$$

and

$$
\mathbb{I}_{\left\{\sum_{j=1}^{p} m_{j, t}, \sum_{j=1}^{p} m_{j, t}+1, \ldots\right\}}\left(y_{t}\right) \times \mathbb{I}_{\left\{0, \ldots, y_{t-i}\right\}}\left(m_{i, t}\right)=\mathbb{I}_{\left\{0,1, \ldots, \min \left\{y_{t}-\sum_{j \neq i} m_{j, t}, y_{t-i}\right\}\right\}}\left(m_{i, t}\right),
$$

we recover the original likelihood of the generalized $\operatorname{INAR}(p)$, showing that the introduction of the latent maturations $M_{i, t}$ with the specified distributions is a valid data augmentation scheme (see [8, 9$]$ for a general discussion of data augmentation techniques).

In the next section, we review the needed definitions and properties of the Pitman-Yor process.

## 3 Pitman-Yor process

Let the random probability measure $\mathbb{G} \sim \mathrm{DP}\left(\tau, G_{0}\right)$ be a Dirichlet process [10, 11, 12] with concentration parameter $\tau$ and base measure $G_{0}$. If the random variables $X_{1}, \ldots, X_{n}$, given $\mathbb{G}=G$, are conditionally independent and identically distributed as $G$, then it follows that

$$
\operatorname{Pr}\left\{X_{n+1} \in B \mid X_{1}=x_{1}, \ldots, X_{n}=x_{n}\right\}=\frac{\tau}{\tau+n} G_{0}(B)+\frac{1}{\tau+n} \sum_{i=1}^{n} I_{B}\left(x_{i}\right),
$$

for every Borel set $B$. If we imagine the sequential generation of the $X_{i}$ 's, for $i=1, \ldots, n$, the former expression shows that a value is generated anew from $G_{0}$ with probability proportional to $\tau$, or we repeat one the previously generated values with probability proportional to its multiplicity. Therefore, almost surely, realizations of a Dirichlet process are discrete probability measures, perhaps with denumerable infinite support, depending on the nature of $G_{0}$. Also, this data generating process, known as the Pólya-BlackwellMacQueen urn, implies that the $X_{i}$ 's are "softly clustered", in the sense that in one realization of the process the elements of a subset of the $X_{i}$ 's may have exactly the same value.

The Pitman-Yor process [6] is a generalization of the Dirichlet process which results in a model with added flexibility. Essentially, the Pitman-Yor process modifies the expression of the probability associated with the Pólya-Blackwell-MacQueen urn introducing a new parameter so that the posterior predictive probability becomes

$$
\operatorname{Pr}\left\{X_{n+1} \in B \mid X_{1}=x_{1}, \ldots, X_{n}=x_{n}\right\}=\frac{\tau+k \sigma}{\tau+n} G_{0}(B)+\frac{1}{\tau+n} \sum_{i=1}^{n}\left(1-\frac{\sigma}{\nu_{i}}\right) I_{B}\left(x_{i}\right),
$$

in which $0 \leq \sigma<1$ is the discount parameter, $\tau>-\sigma, k$ is the number of distinct elements in $\left\{X_{1}, \ldots, X_{n}\right\}$, and $\nu_{i}$ is the number of elements in $\left\{X_{1}, \ldots, X_{n}\right\}$ which are equal to $X_{i}$, for $i=1, \ldots, n$. It is well known that $\mathrm{E}[\mathbb{G}(B)]=G_{0}(B)$ and

$$
\operatorname{Var}[\mathbb{G}(B)]=\left(\frac{1-\sigma}{\tau+1}\right) G_{0}(B)\left(1-G_{0}(B)\right),
$$

for every Borel set $B$. Hence, $\mathbb{G}$ is centered on the base probability measure $G_{0}$, while $\tau$ and $\sigma$ control the concentration of $\mathbb{G}$ around $G_{0}$. We use the notation $\mathbb{G} \sim \operatorname{PY}\left(\tau, \sigma, G_{0}\right)$. When $\sigma=0$, we recover the Dirichlet process as a special case.

Pitman [6] derived the distribution of the number of clusters $K$ (the number of distinct $X_{i}$ 's), conditionally on both the concentration parameter $\tau$ and the discount parameter $\sigma$, as

$$
\operatorname{Pr}\{K=k \mid \tau, \sigma\}=\frac{\prod_{i=1}^{k-1}(\tau+i \sigma)}{\sigma^{k} \times(\tau+1)_{n-1}} \times \mathscr{C}(n, k ; \sigma),
$$

in which $(x)_{n}=\Gamma(x+n) / \Gamma(x)$ is the rising factorial and $\mathscr{C}(n, k ; \sigma)$ is the generalized factorial coefficient [13].

In the next section, we use a Pitman-Yor process to model the distribution of the innovation rates in the generalized $\operatorname{INAR}(p)$ model.

## 4 PY-INAR $(p)$ model

The PY-INAR $(p)$ model is as a hierarchical extension of the generalized $\operatorname{INAR}(p)$ model defined in Section 2. Given a random measure $\mathbb{G} \sim \operatorname{PY}\left(\tau, \sigma, G_{0}\right)$, in which $G_{0}$ is a Gamma $\left(a_{0}, b_{0}\right)$ distribution, let the innovation rates $\lambda_{p+1}, \ldots, \lambda_{T}$ be conditionally independent and identically distributed with distribution $\operatorname{Pr}\left\{\lambda_{t} \in B \mid \mathbb{G}=G\right\}=G(B)$.

To complete the $\operatorname{PY}-\operatorname{INAR}(p)$ model, we need to specify the form of the prior distribution for the vector of thinning parameters $\alpha=\left(\alpha_{1}, \ldots, \alpha_{p}\right)$. By comparison with standard results from the theory of the $\operatorname{AR}(p)$ model [14], Du and Li [4] found that in the $\operatorname{INAR}(p)$ model the constraint $\sum_{i=1}^{p} \alpha_{i}<1$ must be fulfilled to guarantee the nonexplosiveness of the process. In their Bayesian analysis of the $\operatorname{INAR}(p)$ model, Neal and Kypraios [5] considered independent beta distributions for the $\alpha_{i}$ 's. Unfortunately, this choice is problematic. For example, in the particular case when the $\alpha_{i}$ 's have independent uniform distributions, it is possible to show that $\operatorname{Pr}\left\{\sum_{i=1}^{p} \alpha_{i}<1\right\}=1 / p$ !, implying that we would be concentrating most of the prior mass on the explosive region even for moderate values of the model order $p$. We circumvent this problem using a prior distribution for $\alpha$ that places all of its mass on the nonexplosive region and still allows us to derive the full conditional distributions of the $\alpha_{i}$ 's in simple closed form. Specifically, we take the prior distribution of $\alpha$ to be a Dirichlet distribution with hyperparameters $\left(a_{1}, \ldots, a_{p} ; a_{p+1}\right)$, and corresponding density

$$
\pi(\alpha)=\frac{\Gamma\left(\sum_{i=1}^{p+1} a_{i}\right)}{\prod_{i=1}^{p+1} \Gamma\left(a_{i}\right)} \prod_{i=1}^{p+1} \alpha_{i}^{a_{i}-1}
$$

in which $a_{i}>0$, for $i=1, \ldots, p+1$, and $\alpha_{p+1}=1-\sum_{i=1}^{p} \alpha_{i}$.

Let $m=\left\{m_{i, t}: i=1, \ldots, p, t=p+1, \ldots, T\right\}$ denote the set of all maturations, and let $\mu_{\mathbb{G}}$ be the distribution of $\mathbb{G}$. Our strategy to derive the full conditionals distributions of the model parameters and latent variables is to consider the marginal distribution

$$
\begin{aligned}
p(y, m, \alpha, \lambda)= & \int p(y, m, \alpha, \lambda \mid G) d \mu_{\mathbb{G}}(G) \\
= & \left\{\prod_{t=p+1}^{T} p\left(y_{t} \mid m_{1, t}, \ldots, m_{p, t}, \lambda_{t}\right) \prod_{i=1}^{p} p\left(m_{i, t} \mid y_{t-i}, \alpha_{i}\right)\right\} \\
& \times \pi(\alpha) \times \int \prod_{t=p+1}^{T} p\left(\lambda_{t} \mid G\right) d \mu_{\mathbb{G}}(G) .
\end{aligned}
$$

From this expression, using the results in Section 3, the derivation of the full conditional distributions is straightforward. In the following expressions, the symbol $\propto$ denotes proportionality up to a suitable normalization factor, and the label "all others" designate the observed counts $y$ and all the other latent variables and model parameters, with the exception of the one under consideration.

Let $\lambda_{\backslash t}$ denote the set $\left\{\lambda_{p+1}, \ldots, \lambda_{T}\right\}$ with the element $\lambda_{t}$ removed. Then, for $t=$ $p+1, \ldots, T$, we have

$$
\lambda_{t} \mid \text { all others } \sim w_{t} \times \operatorname{Gamma}\left(y_{t}-m_{t}+a_{0}, b_{0}+1\right)+\sum_{r \neq t}\left(1-\frac{\sigma}{\nu_{r}}\right) \lambda_{r}^{y_{t}-m_{t}} e^{-\lambda_{r}} \delta_{\left\{\lambda_{r}\right\}},
$$

in which the weight

$$
w_{t}=\frac{\left(\tau+k_{\backslash t} \sigma\right) \times b_{0}^{a_{0}} \times \Gamma\left(y_{t}-m_{t}+a_{0}\right)}{\Gamma\left(a_{0}\right) \times\left(b_{0}+1\right)^{y_{t}-m_{t}+a_{0}}}
$$

$\nu_{r}$ is the number of elements in $\lambda_{\backslash t}$ which are equal to $\lambda_{r}$, and $k_{\backslash t}$ is the number of distinct elements in $\lambda_{\backslash t}$. In this mixture, we suppressed the normalization constant that makes all weights add up to one.

Making the choice $a_{p+1}=1$, we have

$$
\alpha_{i} \mid \text { all others } \sim \operatorname{TBeta}\left(a_{i}+\sum_{t=p+1}^{T} m_{i, t}, 1+\sum_{t=p+1}^{T}\left(y_{t-i}-m_{i, t}\right), 1-\sum_{j \neq i} \alpha_{j}\right),
$$

for $i=1, \ldots, p$, in which TBeta denotes the right truncated Beta distribution with support ( $0,1-\sum_{j \neq i}^{p} \alpha_{j}$ ).

For the latent maturations, we find

$$
\begin{aligned}
p\left(m_{i, t} \mid \text { all others }\right) \propto & \frac{1}{\left(m_{i, t}\right)!\left(y_{t}-\sum_{j=1}^{p} m_{j, t}\right)!\left(y_{t-i}-m_{i, t}\right)!}\left(\frac{\alpha_{i}}{\lambda_{t}\left(1-\alpha_{i}\right)}\right)^{m_{i, t}} \\
& \times \mathbb{I}_{\left\{0,1, \ldots, \min \left\{y_{t}-\sum_{j \neq i} m_{j, t}, y_{t-i}\right\}\right\}}\left(m_{i, t}\right) .
\end{aligned}
$$

To explore the posterior distribution of the model, we build a Gibbs sampler [15] using these full conditional distributions. Escobar and West [16] showed, in a similar context,
that we can improve mixing by resampling simultaneously the values of all $\lambda_{t}$ 's inside the same cluster at the end of each iteration of the Gibbs sampler. Letting $\left(\lambda_{1}^{*}, \ldots, \lambda_{k}^{*}\right)$ be the $k$ unique values among $\left(\lambda_{p+1}, \ldots, \lambda_{T}\right)$, define the number of occupants of cluster $j$ by $n_{j}=\sum_{t=p+1}^{T} \mathbb{I}_{\left\{\lambda_{j}^{*}\right\}}\left(\lambda_{t}\right)$, for $j=1, \ldots, k$. It follows that

$$
\lambda_{j}^{*} \mid \text { all others } \sim \operatorname{Gamma}\left(a_{0}+\sum_{t=p+1}^{T}\left(y_{t}-\sum_{j=1}^{p} m_{j, t}\right) \cdot \mathbb{I}_{\left\{\lambda_{j}^{*}\right\}}\left(\lambda_{t}\right), b_{0}+n_{j}\right) .
$$

for $j=1, \ldots, k$. At the end of each iteration of the Gibbs sampler, we update the values of all $\lambda_{t}$ 's inside each cluster by the corresponding $\lambda_{j}^{*}$ using this distribution.

## 5 Prior sensitivity

As it is often the case for Bayesian models with nonparametric components, a choice of the prior parameters of the $\mathrm{PY}-\operatorname{INAR}(p)$ model which yields robustness of the posterior distribution is nontrivial [17].

The first aspect to be considered is the fact that the base measure $G_{0}$ plays a crucial role in the determination of the posterior distribution of the number of clusters $K$. This can be seen directly by inspecting the form of the full conditional distributions derived in Section 4 Recalling that $G_{0}$ is a gamma distribution with mean $a_{0} / b_{0}$ and variance $a_{0} / b_{0}^{2}$, from the full conditional distribution of $\lambda_{t}$ one may note that the probability of generating, on each iteration of the Gibbs sampler, a value for $\lambda_{t}$ anew from $G_{0}$ is proportional to

$$
\frac{\left(\tau+k_{\backslash t} \sigma\right) \times b_{0}^{a_{0}} \times \Gamma\left(y_{t}-m_{t}+a_{0}\right)}{\Gamma\left(a_{0}\right)\left(b_{0}+1\right)^{y_{t}-m_{t}+a_{0}}} .
$$

Therefore, supposing that all the other terms are fixed, if we concentrate the mass of $G_{0}$ around zero by making $b_{0} \rightarrow \infty$, this probability decreases to zero. This is not problematic, because it is hardly the case that we want to make such a drastic choice for $G_{0}$. The behavior in the other direction is more revealing, since taking $b_{0} \downarrow 0$, in order to spread the mass of $G_{0}$, also makes the limit of this probability to be zero. Due to this behavior, we need to establish a criterion to choose the hyperparameters of the base measure which avoids these extreme cases.

In our analysis, it is convenient to have a single hyperparameter regulating how the mass of $G_{0}$ is spread over its support. For a given $\lambda_{\max }>0$, we find numerically the values of $a_{0}$ and $b_{0}$ which minimize the Kullback-Leibler divergence between $G_{0}$ and a uniform distribution on the interval $\left[0, \lambda_{\max }\right]$. This Kullback-Leibler divergence can be computed explicitly as

$$
-\log \lambda_{\max }-a_{0} \log b_{0}+\log \Gamma\left(a_{0}\right)-\left(a_{0}-1\right)\left(\log \lambda_{\max }-1\right)+\frac{b_{0} \lambda_{\max }}{2}
$$

In this new parameterization, our goal is to make a sensible choice for $\lambda_{\max }$. It is worth emphasizing that by this procedure we are not truncating the support of $G_{0}$, but only


Figure 1: Formation of the elbows for $\sigma=0.5$ (left) and $\sigma=0.75$ (right). The red dotted lines indicate the chosen values of $\lambda_{\max }$.
using the uniform distribution on the interval $\left[0, \lambda_{\max }\right]$ as a reference for our choice of the base measure hyperparameters $a_{0}$ and $b_{0}$.

Our proposal to choose $\lambda_{\max }$ goes as follows. We fix some value $0 \leq \sigma<1$ for the discount parameter and choose an integer $k_{0}$ as the prior expectation of the number of clusters $K$, which, using the results at the end of Section 3, can be computed explicitly as

$$
\mathrm{E}[K]= \begin{cases}\tau \times(\psi(\tau+T-p)-\psi(\tau)) & \text { if } \quad \sigma=0 \\ \left((\tau+\sigma)_{T-p} /\left(\sigma \times(\tau+1)_{T-p-1}\right)\right)-\tau / \sigma & \text { if } \quad \sigma>0\end{cases}
$$

in which $\psi(x)$ is the digamma function (see [6] for a derivation of this result). Next, we find the value of the concentration parameter $\tau$ by solving $\mathrm{E}[K]=k_{0}$ numerically. After this, for each $\lambda_{\max }$ in a grid of values, we run the Gibbs sampler and compute the posterior expectation of the number of clusters $\mathrm{E}[K \mid y]$. Finally, in the corresponding graph, we look for the value of $\lambda_{\max }$ located at the "elbow" of the curve, that is, the value of $\lambda_{\max }$ at which the values of $\mathrm{E}[K \mid y]$ level off.

As an explicit example of this graphical criterion in action, we used the functional form of a first-order model with thinning parameter $\alpha=0.15$ to simulate a time series of length $T=1000$, for which the distribution of the innovations is a symmetric mixture of three Poisson distributions with parameters 1,8 , and 15 . Figure 1 shows the formations of the elbows for two values of the discount parameter: $\sigma=0.5$ and $\sigma=0.75$.

For the simulated time series, Figures 2, 3, 4, and 5 display the behavior of the posterior distributions obtained using the elbow method for $\left(k_{0}, \sigma\right) \in\{4,10,16,30\} \times$ $\{0,0.25,0.5,0.75\}$. These figures make the relation between the choice of the value of the discount parameter $\sigma$ and the achieved robustness of the posterior distribution quite explicit: as we increase the value of the discount parameter $\sigma$, the posterior becomes insensitive to the choice of $k_{0}$. In particular, for $\sigma=0.75$, the posterior mode is always

|  | Simulated |  |  |
| :---: | :---: | :---: | :---: |
| Predicted | 1 | 2 | 3 |
| 1 | 297 | 32 | 0 |
| 2 | 11 | 217 | 42 |
| 3 | 0 | 84 | 316 |

Table 1: Confusion matrix for the cluster assignments.
near 3, which is the number of components used in the distribution of the innovations of the simulated time series.

Once we understand the influence of the prior parameters on the robustness of the posterior distribution, an interesting question is how to get a point estimate for the distribution of clusters, in the sense that each $\lambda_{t}$, for $t=p+1, \ldots, T$, would be assigned to one of the available clusters.

From the Gibbs sampler, it is easy to obtain an approximation for the probabilities $d_{r t}=\operatorname{Pr}\left\{\lambda_{r} \neq \lambda_{t} \mid y\right\}$, for $r, t=p+1, \ldots, T$. These probabilities define a dissimilarity matrix $D=\left(d_{r t}\right)$ among the innovation rates. Although $D$ is not a distance matrix, we can use it as a starting point to represent the innovation rates in a two dimensional Euclidean space using the technique of metric multidimensional scaling (see [18] for a general discussion). From this two dimensional representation, we build a dendrogram, which is appropriately cut in order to define three clusters, allowing us to assign a single cluster label to each innovation rate. Table 1 displays the confusion matrix of this assignment, showing that $83 \%$ of the innovations were grouped correctly in the clusters which correspond to the mixture components used to simulate the time series. A visual representation of this difference between assigned and simulated clusters is given on Figure 6

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Figure 2: Posterior distributions of the number of clusters $K$ for the simulated time series with $\sigma=0$ and $k_{0}=4,10,16,30$. The red dotted lines indicate the value of $k_{0}$.


Figure 3: Posterior distributions of the number of clusters $K$ for the simulated time series with $\sigma=0.25$ and $k_{0}=4,10,16,30$. The red dotted lines indicate the value of $k_{0}$.


Figure 4: Posterior distributions of the number of clusters $K$ for the simulated time series with $\sigma=0.5$ and $k_{0}=4,10,16,30$. The red dotted lines indicate the value of $k_{0}$.


Figure 5: Posterior distributions of the number of clusters $K$ for the simulated time series with $\sigma=0.75$ and $k_{0}=4,10,16,30$. The red dotted lines indicate the value of $k_{0}$.


Figure 6: Swallow tails obtained by multidimensional scaling. The left and right figures show the assigned and simulated cluster labels, respectively

