# Convergence rate bounds for iterative random functions using one-shot coupling 

Sabrina Sixta ${ }^{1}$ (D) Jeffrey S. Rosenthal ${ }^{1}$

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#### Abstract

One-shot coupling is a method of bounding the convergence rate between two copies of a Markov chain in total variation distance, which was first introduced in Roberts and Rosenthal (Process Appl 99:195-208, 2002) and generalized in Madras and Sezer (Bernoulli 16:882-908, 2010). The method is divided into two parts: the contraction phase, when the chains converge in expected distance and the coalescing phase, which occurs at the last iteration, when there is an attempt to couple. One-shot coupling does not require the use of any exogenous variables like a drift function or a minorization constant. In this paper, we summarize the one-shot coupling method into the One-Shot Coupling Theorem. We then apply the theorem to two families of Markov chains: the random functional autoregressive process and the autoregressive conditional heteroscedastic process. We provide multiple examples of how the theorem can be used on various models including ones in high dimensions. These examples illustrate how the theorem's conditions can be verified in a straightforward way. The one-shot coupling method appears to generate tight geometric convergence rate bounds.


Keywords One-shot coupling • Convergence rate • Iterated random functions • Markov chain • Total variation distance • Gibbs sampler

## 1 Introduction

The study of Markov chain convergence rates focuses on evaluating how fast a positive recurrent Markov chain converges to its stationary distribution. On one hand, a great deal of progress has been made in bounding the convergence rate for Markov chains defined in discrete state spaces (Saloff-Coste 1997; Rosenthal 1995a, 2016). On the other hand, despite the major developments made in bounding Markov chains in continuous state space, many applications of continuous state space Markov chains do not have established convergence rate bounds. For example, convergence rate bounds applied to Markov chain Monte Carlo (MCMC) models, which are useful for deciding the size of the burnin period (Hobert and Jones 2001; Geyer 2011), do not have known upper bounds on their convergence rate (Geyer 2011).

[^0]Users need to rely on ad-hoc convergence diagnostics (e.g., Gelman and Rubin 1992), which offer no guarantees.

Methods using the drift and minorization conditions (e.g., Rosenthal 1995b; Baxendale 2005), which guarantee geometric ergodicity (Definition 2.5), are the most studied techniques for bounding Markov chains in continuous state space (Roberts and Rosenthal 2004; Hobert and Jones 2001). The minorization condition is satisfied for a Markov chain $\left\{X_{n}\right\}_{n \geq 1}$ under the following circumstances: there exists a small set $K$, a probability measure $Q$ and a positive number $\epsilon>0$ such that $P\left(\cdot \mid X_{n}=x\right) \geq \epsilon Q(\cdot)$ for $x \in K$. The drift condition is satisfied if there exists a positive function $V$, and constants $\alpha>1$ such that $E\left[V\left(X_{n+1}\right) \mid X_{n}=x\right] \leq V(x) / \alpha$ (Meyn and Tweedie 1993; Roberts and Rosenthal 2004). Bounds generated using the drift and minorization conditions have been applied to a wide array of problems such as (Rosenthal 1996; Tan et al. 2013; Hobert and Jones 2001).

Despite the widespread use of bounds generated by the drift and minorization conditions, there are drawbacks. First, it can be a challenge to identify a small set $K$ and drift function $V$ (Madras and Sezer 2010). Second, it is shown in Qin and Hobert (2020) based on results from Jerison (2016) that
bounds that use the minorization condition do not scale well in high dimensions.

Alternatively, methods for finding Markov chain convergence rate bounds on the Wasserstein distance have been shown to scale well in high dimensions (Durmus and Moulines 2015; Qin and Hobert 2020), so bounding the total variation distance by first bounding the Wasserstein distance is a common technique used in the literature (Qin and Hobert 2020; Madras and Sezer 2010; Jin and Hobert 2021).

One-shot coupling, which bounds the Wasserstein distance as an intermediate step (Madras and Sezer 2010), provides an upper bound on the convergence rate in total variation distance of a Markov chain. This method does not need to identify any exogenous sets or functions, like the drift and minorization conditions. Further, the one-shot coupling method has already been shown to scale well in certain highdimensional examples Roberts and Rosenthal (2002); Pillai and Smith (2017) and will be shown in this paper to scale well in high dimensions for the Bayesian regression Gibbs sampler (Example 4.7) and the Bayesian location Gibbs sampler (Example 4.13).

The one-shot coupling method described in Roberts and Rosenthal (2002) works by first converging the expected distance between two copies of a Markov chain. At the last iteration, the probability of coupling is evaluated when the expected distance between the copies is small. This contrasts with the drift and minorization technique, which attempts to couple the two Markov chain copies every time they enter some fixed small set $K$.

In this paper, we introduce the One-Shot Coupling Theorem 3.1, which aims to summarize the method defined in Roberts and Rosenthal (2002) and Madras and Sezer (2010) for straightforward applications. The One-Shot Coupling Theorem is used as the foundation for bounding the convergence rate for all of the examples in this paper, which can be partitioned into two families: the random functional autoregressive process and the ARCH process. In Sect. 4, we introduce the Sideways Theorem 4.2, which is new and is an application of the One-Shot Coupling Theorem. We apply it to various examples of random functional autoregressive processes (Definition 4.1). In Sect. 5, we provide convergence rate bounds using the One-Shot Coupling Theorem to various ARCH processes (Definition 5.1).

Proofs for the theorems presented in this paper are found in the appendix, Sect. 5.2. The code used to generate all of the tables and calculations can be found on github.com/sixter/OneShotCoupling.

## 2 Background and notation

Let $\left\{X_{n}\right\}_{n \geq 1}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1}$ be two copies of the Markov chain over the state space $\mathcal{X}$ and define $\mathcal{L}\left(X_{n}\right)$ to be the distribution
of the random variable $X_{n}$. We define $\pi$ to be the stationary distribution of the Markov chain.

### 2.1 Total variation distance

We are interested in measuring the distance between the distribution of two Markov chains. To measure this, we use the total variation metric.

Definition 2.1 (Total variation distance) The total variation distance between the laws of two random variables, $X$ and $X^{\prime}$, defined on the state space $\mathcal{X}$ is

$$
\left\|\mathcal{L}(X)-\mathcal{L}\left(X^{\prime}\right)\right\|=\sup _{A \subseteq \mathcal{X}}\left|P(X \in A)-P\left(X^{\prime} \in A\right)\right|
$$

where $\mathcal{L}(X)$ represents the distribution of the random variable $X$ and $A$ is a measurable set.

If the random variables, $X, X^{\prime} \in \mathbb{R}$ have defined density functions $f_{X}, f_{X^{\prime}}$ over the reference measure $\lambda$,
$\left\|\mathcal{L}(X)-\mathcal{L}\left(X^{\prime}\right)\right\|=\frac{1}{2} \int_{\mathbb{R}}\left|f_{X}(x)-f_{X^{\prime}}(x)\right| \lambda(d x)$
Total variation distance has natural probability interpretations. It is the maximum difference in probabilities of an event. It is the error in an expected bounded loss function when a given measure is used as a proxy for another (Gibbs and Su 2002). Finally, it can be seen as the percentage of samples of $\mathcal{L}(X)$ which cannot be regarded as samples from $\mathcal{L}\left(X^{\prime}\right)$ (Proposition 3(g) Roberts and Rosenthal 2004).

Historically, total variation distance was the common metric for measuring Markov chain convergence rates (Roberts and Rosenthal 2004; Meyn and Tweedie 1993; Jones Jones2004; Hobert and Jones 2001), and hence, there is a rich literature of attributes that can be deduced from finding convergence rates in total variation (Jin and Tan 2020). For example, mixing times in total variation distance can be used to determine whether the Markov chain is asymptotically uncorrelated (Theorem 2 of Jones (Jones2004)), uniformly integrable (Theorem 3 of Jones (Jones2004)), whether the central limit theorem (CLT) applies (Theorem 9 of Jones (Jones2004) or Section 5.2 of Roberts and Rosenthal (2004)), or whether it is convergent based on the total variation mixing times of another Markov chain (Theorem 8 of Dyer et al. (2006)).

The following are properties of total variation, which will be used in conjunction with the One-Shot Coupling Theorem 3.1 to establish upper bounds on the convergence rate for the examples in this paper.

Proposition 2.2 states that the total variation between two random variables is equal to the total variation of any invertible transform of the same random variables. This proposition
resembles Lemma 4.13 of Levin et al. (2017) and Lemma 3 of Madras and Sezer (2010).

Proposition 2.2 Let $X, X^{\prime} \in \mathcal{X}$ be two random variables and let $g: \mathcal{X} \rightarrow \mathcal{Y}$ be an invertible and measurable function. Then,

$$
\begin{equation*}
\left\|\mathcal{L}(g(X))-\mathcal{L}\left(g\left(X^{\prime}\right)\right)\right\|=\left\|\mathcal{L}(X)-\mathcal{L}\left(X^{\prime}\right)\right\| \tag{2}
\end{equation*}
$$

The proof is in Sect. 1.
In general, for a measurable (not necessarily invertible) function $g, g^{-1}(f(\mathcal{B})) \subset \mathcal{B}$, so the third equality in the proof becomes $\leq$ and
$\left\|\mathcal{L}(g(X))-\mathcal{L}\left(g\left(X^{\prime}\right)\right)\right\| \leq\left\|\mathcal{L}(X)-\mathcal{L}\left(X^{\prime}\right)\right\|$
Proposition 2.3 states that the total variation distance between two random variables is bounded above by the expected value of the conditional random variable.

Proposition 2.3 Let $X, X^{\prime}$ be two random variables with corresponding $\sigma$-field $\mathcal{B}$ and $Y \in \mathcal{Y}$ be some related random variable. Then

$$
\left\|\mathcal{L}(X)-\mathcal{L}\left(X^{\prime}\right)\right\| \leq E\left[\left\|\mathcal{L}(X \mid Y)-\mathcal{L}\left(X^{\prime} \mid Y\right)\right\|\right]
$$

## The proof is in Section A.

Proposition 2.4 states that the convergence rate of a Markov chain in $\mathbb{R}^{d}$ with independent coordinates is $d$ times the maximum coordinate-wise convergence rate. This proposition is an application of inequality 1.2 of Reiss (1981).

Proposition 2.4 Let $\left\{\vec{X}_{n}\right\}_{n \geq 1} \in \mathbb{R}^{d}$ be a Markov chain such that each coordinate is independent of the other coordinates, $X_{i, n} \Perp X_{j, n}, i \neq j$. Further suppose that for two copies of the Markov chain $\left\{\vec{X}_{n}\right\}_{n \geq 1}$ and $\left\{\vec{X}_{n}^{\prime}\right\}_{n \geq 1}$, $\max _{1 \leq i \leq d}\left\|\mathcal{L}\left(X_{i, n}\right)-\mathcal{L}\left(X_{i, n}^{\prime}\right)\right\| \leq A r^{n}$ for some $A \in \mathbb{R}_{+}$ and $r \in(0,1)$. Then,
$\left\|\mathcal{L}\left(\vec{X}_{n}\right)-\mathcal{L}\left(\vec{X}_{n}^{\prime}\right)\right\| \leq d A r^{n}$

## The proof is in Section A.

In this paper, we establish convergence bounds for Markov chains that are geometrically ergodic in total variation distance.

Definition 2.5 (Geometric ergodicity) Let $\left\{X_{n}\right\}_{n \geq 1}$ be a Markov chain with a stationary distribution $\pi$. The Markov chain is geometrically ergodic if there exists a $\rho<1$ and a function $M(x)<\infty, \pi$-a.e. such that for $X_{0}=x$,

$$
\begin{equation*}
\left\|\mathcal{L}\left(X_{n}\right)-\pi\right\| \leq M(x) \rho^{n} \tag{4}
\end{equation*}
$$

The geometric rate of convergence for $X_{n}$ is defined as $\rho^{*}=$ $\inf \{\rho:$ Eq. 4 holds $\}$.

Proposition 4 of Qin and Hobert (2020) states that for any sequence of drift and minorization conditions, the geometric convergence rate $\rho$ established by the Rosenthal bound (Theorem 12 of Rosenthal (1995b)) will increase at an exponential rate for the autoregressive normal process in $\mathbb{R}^{d}$ as the dimension $d \rightarrow \infty$. This finding suggests that convergence bounds that use the drift and minorization condition do not scale well in dimension (see Lemma 3 and discussion in Qin and Hobert (2020)). However, Proposition 2.4 shows that since each coordinate in this example is independent, the geometric convergence $\rho$ rate is indeed invariant to dimension, regardless of the bounding approach. Thus, a drift and minorization bound, including the Rosenthal bound, can easily be applied to the autoregressive normal process in $\mathbb{R}$ and then extended to $\mathbb{R}^{d}$ using Proposition 2.4. To see Proposition 2.4 applied to the autoregressive normal process in $\mathbb{R}^{d}$, see Example 4.21.

### 2.2 Wasserstein distance

Let $X, X^{\prime} \in \mathbb{R}$ be two random variables equipped with the Euclidean distance. The Wasserstein-1 distance can be defined as follows,

$$
\begin{gathered}
W\left(\mathcal{L}(X), \mathcal{L}\left(X^{\prime}\right)\right)=\inf \left\{E\left[\left|Y-Y^{\prime}\right|\right]: \mathcal{L}(X)=\mathcal{L}(Y)\right. \\
\text { and } \left.\mathcal{L}\left(X^{\prime}\right)=\mathcal{L}\left(Y^{\prime}\right)\right\}
\end{gathered}
$$

In comparison with total variation distance, there is not as much literature dedicated to Markov properties that can be derived from the convergence in Wasserstein distance (Jin and Tan 2020). However, this literature is growing. For example, Jin and Tan provide sufficient conditions in Jin and Tan (2020) for the CLT based on convergence in Wasserstein distance (see Theorems 9 and 10).

### 2.3 Coupling

Total variation can also be defined in terms of the coupling characterization (Gibbs and Su 2002),

$$
\begin{aligned}
\left\|\mathcal{L}(X)-\mathcal{L}\left(X^{\prime}\right)\right\|= & \inf \left\{P\left(Y \neq Y^{\prime}\right) \mid \mathcal{L}(X)=\mathcal{L}(Y)\right. \\
& \text { and } \left.\mathcal{L}\left(X^{\prime}\right) \neq \mathcal{L}\left(Y^{\prime}\right)\right\}
\end{aligned}
$$

The total variation metric measures the distance between two distributions, but is invariant to how these measures are jointly distributed. For example, let $X \sim N(0,1)$ and $X^{\prime} \sim N(1,1)$ be two random variables. Regardless of whether $X$ and $X^{\prime}$ are highly dependent, for example if $X=X^{\prime}+1$ or if $X, X^{\prime}$ are independent, their total variation distance would be the same. The Nummelin splitting
technique makes use of this by constructing alternative random variables, $Y$ and $Y^{\prime}$, such that the marginal distributions are the same $\mathcal{L}(X)=\mathcal{L}(Y), \mathcal{L}\left(X^{\prime}\right)=\mathcal{L}\left(Y^{\prime}\right)$, and the probability that they are unequal is minimized. This technique was first shown in Nummelin (1978). See Rosenthal (1995a) or Meyn and Tweedie (1993) for an explanation. Finally, note that the theory on maximal coupling guarantees that there exists alternative random variables $Y, Y^{\prime}$ as defined above, such that $\left\|\mathcal{L}(X)-\mathcal{L}\left(X^{\prime}\right)\right\|=P\left(Y \neq Y^{\prime}\right)$ Böttcher (2017).

Coupling techniques are widely used to calculate total variation upper bounds on Markov chains (Rosenthal 2016, 1995b, a; Roberts and Rosenthal 2004; Rosenthal 1996; Yang and Rosenthal 2019). Let $\left\{X_{n}\right\}_{n \geq 1}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1}$ be two copies of a Markov chain. If we want to use the coupling characterization for finding an upper bound on the total variation distance, we must also make sure that the faithful coupling condition holds (see Section 2 of Rosenthal (1997)). That is, for any measurable set $A \in \mathcal{X}$,

$$
\begin{aligned}
P\left(X_{n+1} \in A: X_{n}\right. & \left.=x \text { and } X_{n}^{\prime}=x^{\prime}\right) \\
& =P\left(X_{n+1} \in A: X_{n}=x\right) \\
P\left(X_{n+1}^{\prime} \in A: X_{n}\right. & \left.=x \text { and } X_{n}^{\prime}=x^{\prime}\right) \\
& =P\left(X_{n+1}^{\prime} \in A: X_{n}^{\prime}=x^{\prime}\right)
\end{aligned}
$$

If the faithful coupling condition holds, then calculating the probability that two chains are unequal at iteration $n$ can be interpreted as the probability that the two chains have not yet coupled by iteration $n$. This is because once the two Markov chains couple, they can be structured so that they are equal forever and so $P\left(X_{n} \neq X_{n}^{\prime}\right)=P(T \geq n)$ where $T=\min \left\{k: X_{k}=X_{k}^{\prime}\right\}$ (Theorem 1 of Rosenthal (1997)). If a minorization condition holds on the Markov chain, then the faithful coupling condition also holds. For one-shot coupling, we do not need faithful coupling, because we only try to couple the chains at the last iteration.

## 3 One-shot coupling

One-shot coupling is an alternative way of applying coupling methods to bound the total variation of two copies of a Markov chain. To apply one-shot coupling, we define a Markov chain in terms of iterated random functions (Diaconis and Freedman 1999). That is, define a family of random functions $\left\{g_{\vec{\theta}}: \vec{\theta} \in \Theta\right\}$ such that $\vec{\theta}_{n}=\left\{\theta_{1, n}, \ldots, \theta_{d, n}\right\}$ is a random vector and
$X_{n}=g_{\vec{\theta}_{n}}\left(X_{n-1}\right)$

The $n$th iteration of the Markov chain can be written in terms of $X_{0}=x$ as follows,
$X_{n}=\left(g_{\vec{\theta}_{n}} \circ g_{\vec{\theta}_{n-1}} \ldots \circ g_{\vec{\theta}_{1}}\right)(x)=g_{\vec{\theta}_{n}}\left(g_{\vec{\theta}_{n-1}}\left(\ldots g_{\vec{\theta}_{1}}(x) \ldots\right)\right)$

Summarizing Section 3 of Roberts and Rosenthal (2002), to find an upper bound on the total variation distance between $X_{N}$ and $X_{N}^{\prime}=g_{\vec{\theta}_{N}^{\prime}}\left(X_{N-1}^{\prime}\right)$ we do the following.

1. Contraction phase: For $n<N$, set $\vec{\theta}_{n}=\vec{\theta}_{n}^{\prime}$ so that the two chains get 'closer' together.
2. Coalescing phase: For $n=N$, we set all but one coordinate of $\vec{\theta}_{n}$ and $\vec{\theta}_{n}^{\prime}$ to equality and attempt to couple $X_{n}$ and $X_{n}^{\prime}$. That is, specify coordinate $j \in\{1, \ldots, d\}$ and set $\theta_{i, n}=\theta_{i, n}^{\prime}$ for all $i \neq j$. Then we attempt to jointly choose $\theta_{j, n}$ and $\theta_{j, n}^{\prime}$, such that

$$
g_{\left(\theta_{1, n}, \ldots, \theta_{j, n}, \ldots, \theta_{d, n}\right)}\left(X_{n-1}\right)=g_{\left(\theta_{1, n}, \ldots, \theta_{j, n}^{\prime}, \ldots, \theta_{d, n}\right)}\left(X_{n-1}^{\prime}\right)
$$

The method used in the contraction phase is also known as the common random number method and is discussed in detail in Section 2.3.1 of Jacob (2021). The contraction phase can also be used to directly generate upper bounds in Wasserstein distance (Jacob 2021; Qin and Hobert 2021; Gibbs 2004) (it is also used to generate bounds on other types of distances like Monge-Kantorovich or Prokhorov Jacob (2021)).

The one-shot coupling method has been applied over a variety of specific examples, namely, a nested gamma model in Jovanovski and Madras (2013), an image restoration model in Jovanovski (2014), and a random walk on the unit sphere in Pillai and Smith (2017).

The contraction and coalescing phase described above is how the one-shot coupling method was first defined in Roberts and Rosenthal (2002). The following theorem summarizes the above method and serves as a general outline for bounding the total variation distance between two Markov chains. The coalescing condition below does not specify how the two chains will couple, unlike the method described above.

Theorem 3.1 (One-Shot Coupling Theorem) Let $\left\{X_{n}\right\}_{n \geq 1}$, $\left\{X_{n}^{\prime}\right\}_{n \geq 1}$ be two copies of a Markov chain such that $X_{n}=$ $g_{\theta_{n}}\left(X_{n-1}\right)$ and $X_{n}^{\prime}=g_{\theta_{n}^{\prime}}\left(X_{n-1}^{\prime}\right)$, where $\left(\theta_{n}, \theta_{n}^{\prime}\right)_{n \geq 1}$ are independent random variables with respect to $n$ and the marginal distribution of $\theta_{n}, \theta_{n}^{\prime} \sim \mathcal{D}$, for some distribution $\mathcal{D}$. Suppose that the following two conditions hold for some non-negative integer $n_{0}$.

1. Contraction condition: There exists a $D \in(0,1)$ such that for any $n \geq n_{0}$ when $\theta_{n+1}=\theta_{n+1}^{\prime} \sim \mathcal{D}$
$E\left[\left|g_{\theta_{n+1}}\left(X_{n}\right)-g_{\theta_{n+1}}\left(X_{n}^{\prime}\right)\right|\right] \leq D E\left[\left|X_{n}-X_{n}^{\prime}\right|\right]$
2. Coalescing condition: There exists $a C>0$ such that for any $n \geq n_{0}$

$$
\left.\| \mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right) \| \leq C E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]
$$

Then the total variation distance between the two Markov chains at iteration $n \geq n_{0}$ is

$$
\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq C D^{n-n_{0}-1} E\left[\left|X_{n_{0}}-X_{n_{0}}^{\prime}\right|\right]
$$

Proof (Proof of the One-Shot Coupling Theorem 3.1) Fix $n \geq n_{0}$. We are interested in finding an upper bound on $\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\|$. To do so, we first generate alternative random variables, $Y_{n}, Y_{n}^{\prime}$ such that

1. for $0 \leq m \leq n_{0}: Y_{m}=X_{m}, Y_{m}^{\prime}=X_{m}^{\prime}$
2. for $n_{0}<m<n: \theta_{m}=\theta_{m}^{\prime} \sim \mathcal{D}$ and $Y_{m}=$ $g_{\theta_{m}}\left(Y_{m-1}\right), Y_{m}^{\prime}=g_{\theta_{m}}\left(Y_{m-1}^{\prime}\right)$.
3. for $m=n: \theta_{m}, \theta_{m}^{\prime} \sim \mathcal{D}$ with an arbitrary joint distribution and $Y_{m}=g_{\theta_{m}}\left(Y_{m-1}\right), Y_{m}^{\prime}=g_{\theta_{m}^{\prime}}\left(Y_{m-1}^{\prime}\right)$

By construction, $Y_{m} \stackrel{d}{=} X_{m}$ and $Y_{m}^{\prime} \stackrel{d}{=} X_{m}^{\prime}$ for $0 \leq m \leq n$.
Next we find an upper bound on the total variation distance between $Y_{n}$ and $Y_{n}^{\prime}$. By the contraction condition for $n_{0} \leq$ $m<n, E\left[\left|g_{\theta_{m+1}}\left(Y_{m}\right)-g_{\theta_{m+1}}\left(Y_{m}^{\prime}\right)\right|\right] \leq D E\left[\left|Y_{m}-Y_{m}^{\prime}\right|\right]$ and so,

$$
\begin{aligned}
E\left[\left|Y_{n-1}-Y_{n-1}^{\prime}\right|\right] & =E\left[\left|g_{\theta_{n-1}}\left(Y_{n-2}\right)-g_{\theta_{n-1}}\left(Y_{n-2}^{\prime}\right)\right|\right] \\
& \leq D E\left[\left|Y_{n-2}-Y_{n-2}^{\prime}\right|\right] \\
& \leq D^{n-n_{0}-1} E\left[\left|Y_{n_{0}}-Y_{n_{0}}^{\prime}\right|\right]
\end{aligned}
$$

By the coalescing condition,

$$
\begin{aligned}
\left.\| \mathcal{L}\left(Y_{n}\right)-\mathcal{L}\left(Y_{n}^{\prime}\right)\right) \| & \leq C E\left[\left|Y_{n-1}-Y_{n-1}^{\prime}\right|\right] \\
& \leq C D^{n-n_{0}-1} E\left[\left|Y_{n_{0}}-Y_{n_{0}}^{\prime}\right|\right] \\
& =C D^{n-n_{0}-1} E\left[\left|X_{n_{0}}-X_{n_{0}}^{\prime}\right|\right]
\end{aligned}
$$

Finally since $Y_{n} \stackrel{d}{=} X_{n}$ and $Y_{n}^{\prime} \stackrel{d}{=} X_{n}^{\prime}$,

$$
\begin{gathered}
\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\|=\left\|\mathcal{L}\left(Y_{n}\right)-\mathcal{L}\left(Y_{n}^{\prime}\right)\right\| \leq C D^{n-n_{0}-1} \\
E\left[\left|X_{n_{0}}-X_{n_{0}}^{\prime}\right|\right]
\end{gathered}
$$

If $\mathcal{L}\left(X_{n}\right)$ has a density function with respect to $X_{n-1}=x$, denoted $f(x, z)$, then Theorem 3.1 can be proven using Wasserstein distance as an intermediary step with the following lemma.

Lemma 3.2 (Theorem 12 of Madras and Sezer (2010)) If $\frac{1}{2} \int_{\mathcal{X}}\left|f(x, z)-f\left(x^{\prime}, z\right)\right| \lambda(d x) \leq C\left|x-x^{\prime}\right|$ holds, then for
$n \geq 0$
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq C W\left(\mathcal{L}\left(X_{n-1}\right), \mathcal{L}\left(X_{n-1}^{\prime}\right)\right)$

If the contraction condition holds, then for $n \geq n_{0}$, $W\left(\mathcal{L}\left(X_{n-1}\right), \mathcal{L}\left(X_{n-1}^{\prime}\right)\right) \leq E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right] \leq D^{n-n_{0}-1}$ $E\left[\left|X_{n_{0}}-X_{n_{0}}^{\prime}\right|\right]$ and the proof of Theorem 3.1 directly follows.

In most cases, $n_{0}=0$. See the GARCH Example 5.12 for an alternative case, $n_{0}=1$.

In general, the contraction condition can be weakened. Theorem 1.1 of Diaconis and Freedman (1999) provides sufficient conditions to guarantee the existence of $D$ as defined in the above theorem. The conditions in Theorem 1 of Steinsaltz (1999), which are called local contractivity and are weaker, could also replace the contraction condition in the above theorem.

To bound the total variation between a Markov chain, $\left\{X_{n}\right\}_{n \geq 1}$, and the corresponding stationary distribution, $\pi$, we set $X_{0}^{\prime} \sim \pi$. This implies that $X_{n}^{\prime} \sim \pi$ and $\| \mathcal{L}\left(X_{n}\right)-$ $\pi \| \leq C D^{n-n_{0}-1} E_{X_{\infty} \sim \pi}\left[\left|X_{n_{0}}-X_{\infty}\right|\right]$ where $C, D$, and $n_{0}$ are satisfied according to the conditions above.

To find an upper bound on $E_{X_{\infty} \sim \pi}\left[\left|X_{n_{0}}-X_{\infty}\right|\right]$, we use the following Lemma 3.4, which uses a drift condition to bound the expected distance between the stationary distribution of a Markov chain and an initial value.

Definition 3.3 (Drift condition) Let $\left\{X_{n}\right\}_{n \geq 1}$ be a Markov chain on $\mathcal{X}$. A drift condition is satisfied if there exists a function $V: \mathcal{X} \rightarrow \mathbb{R}$ and constants $\lambda \in(0,1)$ and $b<\infty$ such that $E\left[V\left(X_{n}\right) \mid X_{n-1}\right] \leq \lambda V\left(X_{n-1}\right)+b$.

Lemma 3.4 Let $\left\{X_{n}\right\}_{n \geq 1}$ be a Markov chain such that a drift condition 3.3 holds with $V(x)=(x+h)^{2}, h \in \mathbb{R}$. The expected distance between $X_{0}$ and $X_{\infty} \sim \pi$ is bounded above as follows, $E\left[\left|X_{\infty}-X_{0}\right|\right] \leq \sqrt{\frac{b}{1-\lambda}}+E\left[\left|X_{0}+h\right|\right]$.

Proof $E\left[\left|X_{\infty}-X_{0}\right|\right] \leq E\left[\left|X_{\infty}+h\right|\right]+E\left[\left|X_{0}+h\right|\right] \leq$ $\sqrt{\frac{b}{1-\lambda}}+E\left[\left|X_{0}+h\right|\right]$. The last inequality holds by Lemma 3.5.

Lemma 3.5 (Proposition 4.3 (i) of Meyn and Tweedie (1993)) If the drift condition holds, then $E_{\pi}[V(X)] \leq \frac{b}{1-\lambda}$. See Sect. C. 3 for a proof.

See Numerical Example 4.18 for an application of Lemma 3.4.

## 4 Random-functional autoregressive processes

The following section proposes the Sideways Theorem to generate upper bounds on the total variation distance for random-functional autoregressive processes.

Definition 4.1 (Random functional autoregressive processes) The sequence $\left\{X_{n}\right\}_{n \geq 1}$ is a random functional autoregressive process if for $g: \mathbb{R}^{2} \rightarrow \mathbb{R}$
$X_{n}=g\left(\theta_{1, n}, X_{n-1}\right)+\theta_{2, n}$
where $\left(\theta_{1, n}, \theta_{2, n}\right) \in \mathbb{R}^{2}$ are random variables and $\left(\theta_{1, n}, \theta_{2, n}\right)$ $\Perp\left(\theta_{1, m}, \theta_{2, m}\right)$ when $n \neq m$.

Theorem 4.2 (Sideways Theorem) Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ be a random-functional autoregressive. Suppose that,

1. Contraction condition: There exists a $D \in(0,1)$ such that for $n \geq 0$,

$$
E\left[\left|g\left(\theta_{1, n+1}, X_{n}\right)-g\left(\theta_{1, n+1}, X_{n}^{\prime}\right)\right|\right] \leq D E\left[\left|X_{n}-X_{n}^{\prime}\right|\right]
$$

2. Attributes of the conditional density $\theta_{2, n} \mid \theta_{1, n}$ : The conditional density of $\theta_{2, n} \mid \theta_{1, n}$
(a) is bounded above: There exists a $K>0$ such that for all $\left(\theta_{1, n}, \theta_{2, n}\right) \in \mathbb{R}^{2}$, the conditional density function of $\theta_{2, n}$ is bounded above by $K, f_{\theta_{2, n}}\left(\theta_{2, n} \mid \theta_{1, n}\right) \leq K$.
(b) has at most $M$ local extrema points that are at most $L>0$ distance apart: For any $\theta_{1, n}$, there are $M$ local maximas and minimas (local extrema points) within the conditional density. The local extrema points are at most $L$ distance apart.
(c) is continuous for any $\theta_{1, n}$

Then an upper bound on the geometric rate of convergence of the Markov chain is $D$ and the total variation distance between the two copies of the Markov chain, $X_{n}, X_{n}^{\prime}$, is bounded above as follows,

$$
\begin{align*}
& \left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \\
& \quad \leq\left(\frac{K(M+1)}{2}+\frac{I_{M>1}}{L}\right) D^{n-1} E\left[\left|X_{0}-X_{0}^{\prime}\right|\right] \tag{5}
\end{align*}
$$

The attributes of the conditional density of $\theta_{2, n} \mid \theta_{1, n}$ serve to prove, by integrating along the $y$-axis or flipping the density sideways, that the coalescing condition is satisfied. To prove the Sideways Theorem, we show that the contraction and coalescing conditions are satisfied and then apply the One-Shot Coupling Theorem 3.1.

Lemma 4.3 (Coalescing condition) If the density of $\theta_{2, n} \mid$ $\theta_{1, n}$ for any $\theta_{1, n}$ is (1) bounded above, (2) has at most $M$
local extrema points that are at most $L$ distance apart, and (3) is continuous then for $n \geq 0$,
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq C E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]$
Where $C=\frac{K(M+1)}{2}+\frac{I_{M>1}}{L}$. See Sect. B.O.1 for a proof.
Proof of Theorem 4.2 If $\left(\theta_{1, n}, \theta_{2, n}\right)=\left(\theta_{1, n}^{\prime}, \theta_{2, n}^{\prime}\right)$ then the contraction condition defined in the One-Shot Coupling Theorem 3.1 holds for $D \in(0,1)$ and $n \geq 0$. The coalescing condition holds by Lemma 4.3, which can be applied when condition 2 is satisfied (attributes of the conditional density of $\theta_{2, n} \mid \theta_{1, n}$ ). By the One-Shot Coupling Theorem 3.1, the total variation distance between two copies of the process can be bounded above using Eq. 5 .

In Guibourg et al. (2012), it is shown that when the function $g$ is deterministic ( $g$ is a function of $X_{n-1}$ only and not $\theta_{1, n}$ ) and given the same assumptions on $\theta_{2, n}$, the upper bound on the geometric rate of convergence is $D$ (see Corollary 8 and Example 9 of Guibourg et al. (2012)). This matches the results from our theorem.

Note that the Sideways Theorem 4.2 provides an upper bound on total variation distance, but does not imply the existence of a stationary distribution for the Markov chain. To develop the intuition for this, first note that convergence in total variation distance implies convergence in distribution (Gibbs and Su 2002 ). Suppose that $\mathcal{L}\left(X_{n}\right), \mathcal{L}\left(X_{n}^{\prime}\right)$ have distribution functions, $F_{n}, F_{n}^{\prime}$, then by Helly's Selection Theorem (see Lemma 11.1.8 of Rosenthal (2016)), a right continuous function $F$ exists such that $F_{n} \rightarrow F$ and $F_{n}^{\prime} \rightarrow F$ pointwise. However, the function $F$ may not necessarily be a distribution function. This is an illustration of why a stationary distribution may not exist.

A simple counter example would be the process $X_{n}=$ $\frac{1}{2} X_{n-1}+n+Z_{n}, Z_{n} \sim N(0,1)$. It is clear how the Sideways Theorem 4.2 could generate a geometric convergence bound over two iterations of the process if $E\left[X_{0}-X_{0}^{\prime}\right]<\infty$, but $X_{n}, X_{n}^{\prime} \rightarrow \infty$ almost surely and so there is no stationary distribution. See (Steinsaltz 1999) for more information on sufficient conditions for stationarity.

### 4.1 An example of a nonlinear autoregressive process

Example 4.4 (Nonlinear autoregressive process) This example is discussed in Section 4 of Qin and Hobert (2021). Let $\left\{X_{n}\right\}_{n \geq 1}$ be a Markov chain such that
$X_{n+1}=\frac{1}{2}\left(X_{n}-\sin X_{n}\right)+Z_{n+1}$
where $\left\{Z_{n}\right\}_{n \geq 1} \sim N(0,1)$ are independent and identically distributed (i.i.d.) random variables. In Qin and Hobert
(2021), it is assumed that $\left\{Z_{n}\right\}_{n \geq 1}$ are i.i.d. random variables with a mean of 0 and a variance of 1 .

For $g(x)=\frac{1}{2}(x-\sin (x))$, the derivative is $g^{\prime}(x)=$ $\frac{1}{2}(1-\cos (x))$ and so $\sup _{x \in \mathbb{R}} g^{\prime}(x)=1$. This cannot be used. Instead, we can find a value for $D$ in terms of the second iteration. That is,
$D^{2}=\sup _{x, y} \frac{E\left[\left|X_{n+2}-X_{n+2}^{\prime}\right| \mid X_{n}=x, X_{n}^{\prime}=y\right]}{|x-y|}$
Lemma 4.5 The value of $D$ as defined above can be written as

Proof If $E\left[\left|\theta_{1, n}\right|\right]<1$, then set $D=E\left[\left|\theta_{1, n}\right|\right]$ and so the contraction condition in Theorem 4.2 holds,

$$
\begin{aligned}
& E\left[\left|g\left(\theta_{1, n+1}, X_{n}\right)-g\left(\theta_{1, n+1}, X_{n}^{\prime}\right)\right|\right] \\
& \quad=E\left[\left|\theta_{1, n+1} X_{n}-\theta_{1, n+1} X_{n}^{\prime}\right|\right] \leq D E\left[\left|X_{n}-X_{n}^{\prime}\right|\right]
\end{aligned}
$$

Since all of the conditions in Theorem 4.2 are satisfied, Eq. 5 holds.

### 4.3 Bayesian regression Gibbs sampler

Example 4.7 (Bayesian regression Gibbs sampler) Suppose
$D^{2}=\sup _{x, y} \frac{\sqrt{2 h(x, y)^{2}-4 \frac{h(x, y) \sin h(x, y) \cos k(x, y)}{e^{1 / 2}}+\sin ^{2} h(x, y)\left(1+\frac{\cos ^{2} k(x, y)-\sin ^{2} k(x, y)}{e^{2}}\right)}}{\sqrt{2}|x-y|}$
where
$h(x, y)=\frac{1}{4}(y-x+\sin x-\sin y)$
$k(x, y)=\frac{1}{4}(x+y-\sin y-\sin x)$
The proof can be found in Sect. C.1.
Using simulation, we can deduce that $D^{2} \approx 0.813^{2}=$ 0.661 , which closely matches the geometric convergence rate found in Qin and Hobert (2021) for the Wasserstein distance of $D=0.814$.

Using the Sideways Theorem 4.2 notation, $K=\frac{1}{\sqrt{2 \pi}}$ and $M=1$. An upper bound on the total variation distance is
$\left\|\mathcal{L}\left(X_{n+1}\right)-\mathcal{L}\left(X_{n+1}^{\prime}\right)\right\| \leq \frac{1}{\sqrt{2 \pi}} E\left[\left|X_{0}-X_{0}^{\prime}\right|\right] 0.661^{\lfloor N / 2\rfloor}$
Thus, if $X_{0}=1$ and $X_{0}^{\prime}=2$, then after 20 iterations, the total variation distance between the two processes will be less than 0.01 .

### 4.2 Random-coefficient autoregressive models

Corollary 4.6 Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ be a random-coefficient autoregressive model. That is, $X_{n}$ is of the following form
$X_{n}=\theta_{1, n} X_{n-1}+\theta_{2, n}$
where $\left(\theta_{1, n}, \theta_{2, n}\right) \Perp\left(\theta_{1, m}, \theta_{2, m}\right)$ when $n \neq m$. If we replace the contraction condition of the Sideways Theorem 4.2 with

1. $E\left[\left|\theta_{1, n}\right|\right]<1$

Then Eq. 5 holds for $D=E\left[\left|\theta_{1, n}\right|\right]$.
we have the following observed data $Y \in \mathbb{R}^{k}$ and $X \in \mathbb{R}^{k \times p}$ where
$Y \mid \beta, \sigma^{2} \sim N_{k}\left(X \beta, \sigma^{2} I_{k}\right)$
for unknown parameters $\beta \in \mathbb{R}^{p}, \sigma^{2} \in \mathbb{R}$. Suppose we apply the prior distributions on the unknown parameters,

- $\beta \left\lvert\, \sigma^{2} \sim N_{p}\left(0_{p}, \frac{\sigma^{2}}{\lambda} I_{p}\right)\right.$, where $\lambda>0$ is known
- $\pi\left(\sigma^{2}\right) \propto 1 / \sigma^{2}$

The Bayesian regression Gibbs sampler is based on the conditional posterior distributions of $\beta_{n}, \sigma_{n}^{2}$ and is defined as follows.

- $\beta_{n} \mid \sigma_{n-1}^{2}, Y \sim N_{p}\left(\tilde{\beta}, \sigma_{n-1}^{2} A^{-1}\right)$
- $\sigma_{n}^{2} \mid \beta_{n}, Y \sim \Gamma^{-1}\left(\frac{k+p}{2}, \frac{1}{2}\left[\left(\beta_{n}-\tilde{\beta}\right)^{T} A\left(\beta_{n}-\tilde{\beta}\right)+C\right]\right)$. $\Gamma^{-1}(\alpha, \beta)$ represents the inverse gamma distribution.
where $A=X^{T} X+\lambda I_{p}$ is positive semi-definite, $\tilde{\beta}=$ $A^{-1} X^{T} Y$, and $C=Y^{T}\left(I_{k}-X A^{-1} X^{T}\right) Y$.

The following theorem gives an upper bound on the convergence rate of the Bayesian regression Gibbs sampler.

Theorem 4.8 For two copies of the Bayesian regression Gibbs sampler, $\left(\beta_{n}, \sigma_{n}\right)$ and $\left(\beta_{n}^{\prime}, \sigma_{n}^{\prime 2}\right)$, defined in Example 4.7,

$$
\begin{align*}
& \left\|\mathcal{L}\left(\beta_{n}, \sigma_{n}\right)-\mathcal{L}\left(\beta_{n}^{\prime}, \sigma_{n}^{\prime 2}\right)\right\| \\
& \quad \leq K E\left[\left|\sigma_{0}^{2}-\sigma_{0}^{\prime 2}\right|\right]\left(\frac{p}{k+p-2}\right)^{n-1} \tag{6}
\end{align*}
$$

where $K=\frac{(C / 2)^{\frac{k+2 p}{2}}}{\Gamma\left(\frac{k+2 p}{2}\right)}\left(\frac{k+2 p+2}{C}\right)^{\frac{k+2 p}{2}+1} e^{-\frac{k+2 p+2}{2}}$.

In Theorem 3.1 of Rajaratnam and Sparks (2015), it was shown than for the equivalent example and some $0<M_{1} \leq$ $M_{2}$, which are not specified,

$$
\begin{aligned}
M_{1}\left(\frac{p}{k+p-2}\right)^{n} & \leq\left\|\mathcal{L}\left(\beta_{n}, \sigma_{n}\right)-\pi\right\| \\
& \leq M_{2}\left(\frac{p}{k+p-2}\right)^{n}
\end{aligned}
$$

This means that the bound derived from the Corollary 4.6 is sharp up to a constant. The primary difference between Theorem 3.1 in Rajaratnam and Sparks (2015) and the bound in Theorem 4.8 is that the latter provides explicit values for the constant, $M_{2}$ and as a result, numerical upper bounds can be calculated.

Before proving Theorem 4.8, we present some lemmas.
Lemma 4.9 The variable $\sigma_{n}^{2}$ can be written as a randomcoefficient autoregressive process, $\sigma_{n}^{2}=X_{n} Y_{n} \sigma_{n-1}^{2}+Y_{n}$ where $X_{n} \sim \Gamma\left(\frac{p}{2}, \frac{C}{2}\right)$ and $Y_{n} \sim \Gamma^{-1}\left(\frac{k+p}{2}, \frac{C}{2}\right)$. And so, $\left\|\mathcal{L}\left(\beta_{n}, \sigma_{n}^{2}\right)-\mathcal{L}\left(\beta_{n}^{\prime}, \sigma_{n}^{\prime 2}\right)\right\| \leq\left\|\mathcal{L}\left(\sigma_{n}^{2}\right)-\mathcal{L}\left(\sigma_{n}^{\prime 2}\right)\right\|$.

The proof can be found in C.2. $\Gamma(\alpha, \beta)$ represents the gamma distribution.

Lemma 4.10 (Contraction condition) The Bayesian regression Gibbs sampler satisfies the contraction condition with $D=\left(\frac{p}{k+p-2}\right)$. The proof can be found in C.2.
Lemma 4.11 (Attributes of the conditional density $\theta_{2, n} \mid$ $\theta_{1, n}$ ) For the Bayesian regression Gibbs sampler, $\theta_{2, n} \mid \theta_{1, n}$ has a continuous density, $M=1$ and $K=$ $\frac{(C / 2)^{\frac{k+2 p}{2}}}{\Gamma\left(\frac{k+2 p}{2}\right)}\left(\frac{k+2 p+2}{C}\right)^{\frac{k+2 p}{2}+1} e^{-\frac{k+2 p+2}{2}}$. The proof can be found in C. 2 .

Given the above lemmas, the proof of Theorem 4.8 is straightforward when the Sideways Theorem is applied.

Proof of Theorem 4.8 Let $n \geq 0$.

$$
\begin{aligned}
& \left\|\mathcal{L}\left(\beta_{n}, \sigma_{n}^{2}\right)-\mathcal{L}\left(\beta_{n}^{\prime}, \sigma_{n}^{\prime 2}\right)\right\| \leq\left\|\mathcal{L}\left(\sigma_{n}^{2}\right)-\mathcal{L}\left(\sigma_{n}^{\prime 2}\right)\right\| \\
& \quad \leq K E\left[\left|\sigma_{0}^{2}-\sigma_{0}^{\prime 2}\right|\right]\left(\frac{p}{k+p-2}\right)^{n-1}
\end{aligned}
$$

where $K$ is defined in Lemma 4.11. Lemma 4.9 implies the first inequality. The second inequality is a result of Corollary 4.6, which is satisfied because of the contraction condition (Lemma 4.10) and the properties of the conditional density $\theta_{2, n} \mid \theta_{1, n}$ (Lemma 4.11).

Numerical Example 4.12 (Application of the Bayesian regression Gibbs sampler) Suppose that we are interested in evaluating the delay in getting a PhD (Y), based on age, age squared, sex and whether the student has a child at home
(X). For more information on this problem, see van de Schoot et al. (2013); Smeets et al. (2019). We want to find the upper bound on the total variation distance for a Bayesian regression Gibbs sampler fitted to this model. In this case, there are 333 observed values $(k=333)$ and 4 covariates $(p=4)$. Using the notation from Theorem 4.8, $K=0.0682$. Further suppose that we are interested in evaluating the upper bound between two copies of the Markov chain $X_{n}, X_{n}^{\prime}$ such that $\sigma_{0}^{2}=1$ and $\sigma_{0}^{\prime 2}=1001$. Then,
$\left\|\mathcal{L}\left(\beta_{n}, \sigma_{n}\right)-\mathcal{L}\left(\beta_{n}^{\prime}, \sigma_{n}^{\prime}\right)\right\| \leq 68.16454(0.0119403)^{n-1}$
After three iterations, the total variation distance between the two chains will be less than 0.01.

### 4.4 Bayesian location model Gibbs sampler

Example 4.13 (Bayesian location model Gibbs sampler) Suppose that we are given data points $Y_{1}, \ldots, Y_{J} \sim N\left(\mu, \tau^{-1}\right)$ where $\mu, \tau^{-1}$ are unknown and $J \geq 3$. Let $\mu, \tau^{-1}$ have flat priors on $\mathbb{R}$ and $\mathbb{R}_{+}$. The Gibbs algorithm is based on the conditional posterior distributions of $\mu, \tau^{-1}$, which are defined as follows.

- $\mu_{n+1}=\bar{y}+Z_{n+1} / \sqrt{J \tau_{n}}$
- $\tau_{n+1}^{-1}=\frac{\frac{S}{2}+\frac{J}{2}\left(\bar{y}-\mu_{n+1}\right)^{2}}{G_{n+1}}$
where $Z_{n} \sim N(0,1)$ and $G_{n} \sim \Gamma\left(\frac{J+2}{2}, 1\right)$ are independent and $S=\sum_{i=1}^{n}\left(y_{i}-\bar{y}\right)^{2}$.

The following theorem gives an upper bound on the convergence rate of the Bayesian location model Gibbs sampler.

Theorem 4.14 For two copies of the Bayesian location model Gibbs sampler Example 4.13,

$$
\begin{align*}
& \left\|\mathcal{L}\left(\mu_{n}, \tau_{n}^{-1}\right)-\mathcal{L}\left(\mu_{n}^{\prime}, \tau_{n}^{\prime-1}\right)\right\| \\
& \quad \leq K E\left[\left|\tau_{0}^{-1}-\tau_{0}^{\prime-1}\right|\right]\left(\frac{1}{J}\right)^{n-1} \tag{8}
\end{align*}
$$

where $K=\frac{(S / 2)^{\frac{J-1}{2}}}{\Gamma\left(\frac{J-1}{2}\right)}\left(\frac{S}{J+1}\right)^{-\frac{J-3}{2}} e^{-\frac{J+1}{2}}$.
This bound compares to the one derived in Section 6 of Roberts and Rosenthal (2002) which states that,

$$
\begin{aligned}
& \left\|\mathcal{L}\left(\mu_{n}, \tau_{n}^{-1}\right)-\mathcal{L}\left(\mu_{n}^{\prime}, \tau_{n}^{\prime-1}\right)\right\| \\
& \quad \leq\left(\frac{J}{2}+1\right) E\left[\left|\tau_{0}^{-1}-\tau_{0}^{\prime-1}\right|\right]\left(\frac{1}{J}\right)^{n}
\end{aligned}
$$

Both bounds return the same geometric rate of convergence. However, the magnitude of constant $K$ is difficult to compare
against $\left(\frac{J}{2}+1\right)$ without knowing $S$. Note that the bound derived from Corollary 4.6 is calculated in a systematic way.

Before proving Theorem 4.14, we present some lemmas.
Lemma 4.15 The variable $\tau_{n}^{-1}$ can be written as a randomcoefficient autoregressive process, $\tau_{n}^{-1}=X_{n} Y_{n} \tau_{n-1}^{-1}+Y_{n}$, where $X_{n} \sim \Gamma\left(\frac{1}{2}, \frac{S}{2}\right)$ and $Y_{n} \sim \Gamma^{-1}\left(\frac{J+2}{2}, \frac{S}{2}\right)$. And so, $\left\|\mathcal{L}\left(\mu_{n}, \tau_{n}^{-1}\right)-\mathcal{L}\left(\mu_{n}^{\prime}, \tau_{n}^{\prime-1}\right)\right\| \leq\left\|\mathcal{L}\left(\tau_{n}^{-1}\right)-\mathcal{L}\left(\tau_{n}^{\prime}-1\right)\right\|$. The proof can be found in C.3.

Lemma 4.16 (Contraction condition) The Bayesian location model Gibbs sampler satisfies the contraction condition with $D=\frac{1}{J}$. The proof can be found in C.3.

Lemma 4.17 (Attributes of the conditional density $\theta_{2, n} \mid$ $\left.\theta_{1, n}\right)$ For the Bayesian location model Gibbs sampler, $\theta_{2, n} \mid$ $\theta_{1, n}$ has a continuous density, $M=1$ and
$K=\frac{(S / 2)^{\frac{J-1}{2}}}{\Gamma\left(\frac{J-1}{2}\right)}\left(\frac{S}{J+1}\right)^{-\frac{J-3}{2}} e^{-\frac{J+1}{2}}$
The proof can be found in C.3.
Given the above lemmas, the proof of Theorem 4.14 is straightforward when the Sideways Theorem is applied.

Proof of Theorem 4.14 Note that

$$
\begin{aligned}
& \left\|\mathcal{L}\left(\mu_{n}, \tau_{n}^{-1}\right)-\mathcal{L}\left(\mu_{n}^{\prime}, \tau_{n}^{-1^{\prime}}\right)\right\| \\
& \quad \leq\left\|\mathcal{L}\left(\tau_{n}^{-1}\right)-\mathcal{L}\left(\tau_{n}^{-1^{\prime}}\right)\right\| \leq K E\left[\left|\tau_{0}^{-1}-\tau_{0}^{-1^{\prime}}\right|\right]\left(\frac{1}{J}\right)^{n-1}
\end{aligned}
$$

where $K$ is defined in Lemma 4.17. The first and second inequality are a result of Lemma 4.15 and Corollary 4.6, respectively. Corollary 4.6 is satisfied because of the contraction condition (Lemma 4.16) and the properties of the conditional density $\theta_{2, n} \mid \theta_{1, n}$ (Lemma 4.17).

Numerical Example 4.18 (Application of Bayesian location model Gibbs sampler) Suppose that we are given the girth in inches of a sample of trees (see the trees dataset in $R$ ), $Y_{1}, \ldots, Y_{31} \sim N\left(\mu, \tau^{-1}\right)$, where $\mu, \tau^{-1}$ are unknown. We want to find the upper bound on the total variation distance for the Gibbs sampler model applied to this problem. In this case, the number of datapoints is $31(J=31)$ and using the notation from Theorem 4.14, $K=13.74027$. Further, suppose that we are interested in evaluating the upper bound between a Markov chain $\left(\mu_{n}, \tau_{n}^{-1}\right)$ with initial value $\tau_{0}^{-1}=1$ and the corresponding stationary Markov chain, which is denoted as $\left(\mu_{\infty}, \tau_{\infty}^{-1}\right)$.

By Lemma 4.19, a drift function exists.
Lemma 4.19 For Numerical Example 4.18, the following drift condition holds,

$$
E\left[\left(\tau_{n}^{-1}+0.5248723\right)^{2} \mid \tau_{n-1}^{-1}\right]
$$

$$
\leq 0.6583702\left(\tau_{n-1}^{-1}+0.5248723\right)^{2}+106.3874
$$

The proof can be found in C.3.
So by Lemma 3.4,
$E\left[\left|\tau_{\infty}^{-1}-\tau_{0}^{-1}\right|\right] \leq 18.12198$
Combining this with Theorem 4.14,

$$
\begin{aligned}
& \left\|\mathcal{L}\left(\mu_{n}, \tau_{n}^{-1}\right)-\mathcal{L}\left(\mu_{\infty}, \tau_{\infty}^{-1}\right)\right\| \\
& \quad \leq 13.74027 \times 18.12198\left(\frac{1}{31}\right)^{n-1} \doteq 249\left(\frac{1}{31}\right)^{n-1}
\end{aligned}
$$

After 4 iterations, the total variation distance between the two chains will be less than 0.01 . This bound compares to the bound derived in Roberts and Rosenthal (2002), which, combined with Eq. 10 , states that $\left\|\mathcal{L}\left(\mu_{n}, \tau_{n}^{-1}\right)-\mathcal{L}\left(\mu_{\infty}, \tau_{\infty}^{-1}\right)\right\| \leq$ $299\left(\frac{1}{31}\right)^{n}$.

### 4.5 Autoregressive normal process

Example 4.20 (Autoregressive normal process in $\mathbb{R}$ ) Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ be an autoregressive normal process. Then for i.i.d. $Z_{n} \sim N(0,1)$,
$X_{n}=\frac{1}{2} X_{n-1}+\sqrt{\frac{3}{4}} Z_{n}$
In this case, $\theta_{1, n}=\frac{1}{2}$ and $\theta_{2, n}=\sqrt{\frac{3}{4}} Z_{n}$. The density of $\theta_{2, n}$ is continuous and uni-modal and $K=\sqrt{\frac{2}{3 \pi}}$. By Corollary 4.6,
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \sqrt{\frac{2}{3 \pi}} E\left[\left|X_{0}-X_{0}^{\prime}\right|\right]\left(\frac{1}{2}\right)^{n-1}$
It is known that the geometric rate of convergence for the autoregressive normal process is $1 / 2$ (Qin and Hobert 2020), so once again the Sideways Theorem 4.2 generates tight geometric convergence rates up to a constant.

When comparing the upper bound with the actual total variation distance, note that if $X_{0}=x_{0}$ is known, $X_{n} \sim$ $N\left(\frac{x_{0}}{2^{n}}, 1-\frac{1}{4^{n}}\right)$. Thus, the total variation distance between two copies of an autoregressive normal process $X_{n}, X_{n}^{\prime}$ where the initial values are known, $X_{0}=x_{0}$ and $X_{0}^{\prime}=x_{0}^{\prime}$, is as follows (see Section 2 of Roberts and Rosenthal (2002)),
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\|=1-2 \Phi\left(-\frac{\left|x_{0}-x_{0}^{\prime}\right|}{2^{n+1} \sqrt{1-\frac{1}{4^{n}}}}\right)$
Figure 1 shows how the upper bound for the autoregressive normal process using Eq. 11 compares to the actual total

Fig. 1 This figure compares the actual value of
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\|$ against the upper bound derived from the Sideways Theorem 4.2, (Eq. 11) when $X_{n}, X_{n}^{\prime}$ are two copies of the autoregressive normal process (i.e.,
$X_{n}=\frac{1}{2} X_{n-1}+\sqrt{\frac{3}{4}} Z_{n}, Z_{n} \sim$ $N(0,1))$ and $x_{0}=0, x_{0}^{\prime}=1$

variation distance when $x_{0}=0$ and $x_{0}^{\prime}=1$. The total variation is less than 0.01 after 6 iteration and the upper bound on the total variation is less than 0.01 after 7 iterations.

In the following section, we extend the above example to higher dimensions.

### 4.6 Processes in $\mathbb{R}^{\text {d }}$

Next we extend the autoregressive normal process as defined above to $\mathbb{R}^{d}$.

To do so, we apply Proposition 2.4 to an autoregressive normal process in $\mathbb{R}^{d}$ with independent coordinates, Example 4.21, and non-independent coordinates, Example 4.22.

Example 4.21 (Autoregressive normal process in $\mathbb{R}^{d}$ with independent coordinates) Let $\left\{\vec{X}_{n}\right\}_{n \geq 1} \in \mathbb{R}^{d}$ be an autoregressive normal process with independent coordinates. Then for i.i.d. $\vec{Z}_{n} \sim N\left(\overrightarrow{0}, I_{d}\right)$,
$\vec{X}_{n}=\frac{1}{2} \vec{X}_{n-1}+\sqrt{\frac{3}{4}} \vec{Z}_{n}$
And if $i \neq j$, then $Z_{i, n} \Perp Z_{j, n}$. Further, $X_{i, n}=\frac{1}{2} X_{i, n-1}+$ $\sqrt{\frac{3}{4}} Z_{i, n}$ for $i \in\{1, \ldots, d\}$ and so by Example 4.20,
$\left\|\mathcal{L}\left(X_{i, n}\right)-\mathcal{L}\left(X_{i, n}^{\prime}\right)\right\| \leq \sqrt{\frac{2}{3 \pi}} E\left[\left|X_{i, 0}-X_{i, 0}^{\prime}\right|\right]\left(\frac{1}{2}\right)^{n-1}$
Since each coordinate is independent and bounded above by the same value, Proposition 2.4 implies that

$$
\left\|\mathcal{L}\left(\vec{X}_{n}\right)-\mathcal{L}\left(\vec{X}_{n}^{\prime}\right)\right\| \leq d \sqrt{\frac{2}{3 \pi}} \sup _{0 \leq i \leq d} E\left[\left|X_{i, 0}-X_{i, 0}^{\prime}\right|\right]\left(\frac{1}{2}\right)^{n-1}
$$

Again, it is known that the geometric rate of convergence for the autoregressive normal process in $\mathbb{R}^{d}$ is $1 / 2$ (Qin and Hobert, 2020).

Finally, to apply numbers to this example, suppose that $\vec{X}_{n}, \vec{X}_{n}^{\prime} \in \mathbb{R}^{100}$ and the initial values of this process are $\vec{X}_{0}=$ $(1, \ldots, 1)$ and $\vec{X}_{0}^{\prime}=(0, \ldots, 0)$. The total variation distance would be bounded above with $\left\|\mathcal{L}\left(\vec{X}_{n+1}\right)-\mathcal{L}\left(\vec{X}_{n+1}^{\prime}\right)\right\| \leq$ $100 \sqrt{\frac{2}{3 \pi}}\left(\frac{1}{2}\right)^{n}$. This means that at 14 iterations the total variation distance would be less than 0.01 .

The following example is a more general version of the above, where $X_{n}$ is a general autoregressive normal process in $\mathbb{R}^{d}$.

Example 4.22 (Autoregressive normal process in $\mathbb{R}^{d}$ ) The random vector $\left\{\vec{X}_{n}\right\}_{n \geq 1} \in \mathbb{R}^{d}$ is an autoregressive normal process if for matrix $A$ and random vector $\vec{W}_{n} \sim N\left(\overrightarrow{0}, \Sigma_{d}^{2}\right)$ ( $\Sigma_{d}^{2}$ is a positive semi-definite matrix)
$\vec{X}_{n}=A \vec{X}_{n-1}+\vec{W}_{n}$

Theorem 4.23 Suppose that $A$ is a diagonalizable matrix. Then for two copies, $\vec{X}_{n}, \vec{X}_{n}^{\prime} \in \mathbb{R}^{d}$, of the autoregressive normal process defined in Example 4.22,

$$
\begin{align*}
& \left\|\mathcal{L}\left(\vec{X}_{n}\right)-\mathcal{L}\left(\vec{X}_{n}^{\prime}\right)\right\| \leq \sqrt{\frac{d}{2 \pi}}\left\|\Sigma_{d}^{-1}\right\|_{2}\|P\|_{2}\left\|P^{-1}\right\|_{2} \\
& \quad E\left[\left\|\vec{X}_{0}-\vec{X}_{0}^{\prime}\right\|_{2}\right] \max _{1 \leq i \leq d}\left|\lambda_{i}\right|^{n} \tag{12}
\end{align*}
$$

where $A=P D P^{-1}$ with $D$ as the corresponding diagonal matrix, $\lambda_{i}$ is the ith eigenvalue of $A$ and $\|\cdot\|_{2}$ denotes the Frobenius norm. The proof can be found in C.4, which uses a modified version of the Sideways Theorem.

Numerical Example 4.24 (Application of the autoregressive normal process in $\mathbb{R}^{d}$ ) To apply numbers to this example, suppose that $\vec{X}_{n}, \vec{X}_{n}^{\prime} \in \mathbb{R}^{100}$ are two copies of the following process $\vec{X}_{n}=A \vec{X}_{n}+\vec{Z}_{n}, \vec{Z}_{n} \sim N(0, A)$ where
$A=\left(\begin{array}{cccccc}\frac{1}{2} & \frac{1}{8} & 0 & \cdots & 0 & 0 \\ \frac{1}{8} & \frac{1}{2} & \frac{1}{8} & \cdots & 0 & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\ 0 & 0 & 0 & \cdots & \frac{1}{8} & \frac{1}{2}\end{array}\right)$
and the initial values of this process are $\vec{X}_{0}=(1, \ldots, 1)$ and $\vec{X}_{0}^{\prime}=(0, \ldots, 0)$. The total variation distance would be bounded above with $\left\|\mathcal{L}\left(\vec{X}_{n}\right)-\mathcal{L}\left(\vec{X}_{n}^{\prime}\right)\right\| \leq 98782.31$ $(0.7498791)^{n}$. This means that after 56 iterations the total variation distance would be less than 0.01 .

## 5 Autoregressive conditional heteroscedastic processes

In this section, we look at bounding the total variation distance between two copies of an ARCH process.

Definition 5.1 (Autoregressive conditional heteroscedastic (ARCH) process) The sequence $\left\{X_{n}\right\}_{n \geq 1}$ is an ARCH process if for $g: \mathbb{R}^{2} \rightarrow \mathbb{R}$
$X_{n}=g\left(\theta_{1, n}, X_{n-1}\right) \theta_{2, n}$
where $\left(\theta_{1, n}, \theta_{2, n}\right) \in \mathbb{R}^{2}$ are random variables and $\left(\theta_{1, n}, \theta_{2, n}\right)$ $\Perp\left(\theta_{1, m}, \theta_{2, m}\right)$ when $n \neq m$.

### 5.1 Application to the LARCH model

Example 5.2 (Linear ARCH process) Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ be a linear ARCH process. Then for i.i.d. $Z_{n}$ and $\beta_{0}, \beta_{1} \in \mathbb{R}$
$X_{n}=\left(\beta_{0}+\beta_{1} X_{n-1}\right) Z_{n}$

See Section 7.3.3 of Doukhan (2018) for more details on this model.

The following theorem provides an upper bound on the convergence rate of two copies of a LARCH process.

Theorem 5.3 Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1} \in \mathbb{R}$ be two copies of the linear ARCH process. Suppose that,

- $\beta_{0}, \beta_{1}>0$ and $Z_{n}>0$ a.s.
- the density of $\log \left(Z_{0}\right)$ is bounded above, has at most $M$ local maxima and minima, and is continuous.

Then, the process is geometrically ergodic if $\beta_{1} E\left[\left|Z_{0}\right|\right]<1$ and an upper bound on the total variation distance between the two processes is,

$$
\begin{align*}
& \left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \frac{\beta_{1}(M+1)}{2 \beta_{0}} \sup _{x} e^{x} f_{Z_{n}}\left(e^{x}\right) \\
& \quad D^{n-1} E\left[\left|X_{0}-X_{0}^{\prime}\right|\right] \tag{14}
\end{align*}
$$

Where $D=\beta_{1} E\left[Z_{0}\right]$
Lemma 7.3.2 of Doukhan (2018) says that if $\beta_{1} E\left[\left|Z_{0}\right|\right]<1$, then a stationary distribution exists. This theorem makes an even stronger assertion that under some additional assumptions, the process will also be geometrically ergodic with geometric convergence rate $D=\beta_{1} E\left[\left|Z_{0}\right|\right]<1$.

Before proving Theorem 5.3, we present some lemmas.
Lemma 5.4 (Contraction condition) The LARCH process satisfies the contraction condition if $D=\beta_{1} E\left[Z_{0}\right]<1$. See Sect. D. 1 for a proof.

Lemma 5.5 (Coalescing condition) Suppose that the density of $\log \left(Z_{0}\right)$ is bounded above, has at most M local maxima and minima and is continuous. Then the LARCH process satisfies the coalescing condition
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq C E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]$
where $n \geq 1$ and $C=\frac{\beta_{1}(M+1)}{2 \beta_{0}} \sup _{x} e^{x} f_{Z_{n}}\left(e^{x}\right)$, See Sect. D. 1 for a proof.

Note that the density of $\log \left(Z_{0}\right)$ is $f_{\log \left(Z_{0}\right)}(x)=e^{x} f_{Z_{0}}\left(e^{x}\right)$.
Proof of Theorem 5.3 Suppose that the assumptions in Theorem 5.3 are satisfied. Then the LARCH model satisfies the contraction condition (Lemma 5.4) and the coalescing condition (Lemma 5.5). By the One-Shot Coupling Theorem 3.1, Eq. 14 holds.

Numerical Example 5.6 We find the convergence rate of Example 10.3.1 of Brockwell and Davis (2002), which is of the form,
$X_{n}^{2}=\left(1+0.5 X_{n-1}^{2}\right) Z_{n}^{2}$
where $Z_{n}^{2} \sim \chi^{2}(1)$. Further let $X_{0}=0.1$ and $X_{0}^{\prime}=1.1$. The density of $\log \left(Z_{n}^{2}\right)$ is $f_{\log \left(Z_{n}^{2}\right)}(x)=(2 \pi)^{-1 / 2} e^{\left(x-e^{x}\right) / 2}$ and so, $\sup f_{\log \left(Z_{n}^{2}\right)}(x)=(2 \pi)^{-1 / 2} e^{\left(0-e^{0}\right) / 2}=\frac{1}{\sqrt{2 \pi e}}$. The density of $\log \left(Z_{n}^{2}\right)$ is also unimodal, so $M=1$. By Theorem 5.3, an upper bound on the total variation distance is
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \frac{1}{\sqrt{8 \pi e}} 0.5^{n-1}$

After three iterations, the total variation distance is less than 0.01. In comparison, Fig. 2 shows how the bound compares to a simulated estimate of the total variation distance for this process.

### 5.2 Application to the asymmetric ARCH model

Example 5.7 (Asymmetric ARCH process) Let $\left\{X_{n}\right\}_{n} n \geq 1$ $\in \mathbb{R}$ be an asymmetric ARCH process. Then for i.i.d. $Z_{n}$
$X_{n}=\sqrt{\left(a X_{n-1}+b\right)^{2}+c^{2}} Z_{n}$
where $a, b, c \in \mathbb{R}$. See Exercise 4.1 of Doukhan (2018) for more details on this process.

The following theorem provides an upper bound on the convergence rate of two copies of an asymmetric ARCH process.
Theorem 5.8 Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1} \in \mathbb{R}$ be two copies of the asymmetric ARCH process defined in Example 5.7. Suppose further that the density of $Z_{n}$ is centred at 0 and is monotonically decreasing around zero (i.e., $\pi(x) \geq \pi(y)$ if $|x|<|y|$.). Then, the process is geometrically ergodic if $|a| E\left[\left|Z_{0}\right|\right]<1$ and an upper bound on the total variation distance between the two processes is
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \frac{|a|}{c} D^{n-1} E\left[\left|X_{0}-X_{0}^{\prime}\right|\right]$
where $D=|a| E\left[\left|Z_{0}\right|\right]$
Exercise 4.1 part 1 of Doukhan (2018) states that the process has a stationary solution if $D=|a| E\left[\left|Z_{0}\right|\right]<1$. Theorem 5.8 shows that under certain additional assumptions on $Z_{n}$ the process will also be geometrically ergodic with a specified quantitative bound.

Before proving Theorem 5.8, we present some lemmas.
Lemma 5.9 (Contraction condition) The asymmetric ARCH process satisfies the contraction condition if $D=|a| E\left[\left|Z_{0}\right|\right]$ $<1$. See Sect. D. 2 for a proof.

Lemma 5.10 (Coalescing condition) Suppose that the density of $Z_{n}$ is centred at 0 and is monotonically decreasing around zero. Then, the asymmetric ARCH process satisfies the coalescing condition

$$
\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq C E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]
$$

where $n \geq 1$ and $C=\frac{|a|}{c}$. See Sect. D. 2 for a proof.
Proof of Theorem 5.8 Suppose that the assumptions in Theorem 5.8 are satisfied. Then the asymmetric ARCH model satisfies the contraction condition (Lemma 5.9) and the coalescing condition (Lemma 5.10). By the One-Shot Coupling Theorem 3.1, Eq. 16 holds.

Numerical Example 5.11 Suppose $a=0.5, b=3, c=5$, $Z_{n} \sim N(0,1)$ and $X_{0}=0, X_{0}^{\prime}=5$. Then by Jensen's inequality, $D=0.5 E\left[\left|Z_{0}\right|\right] \leq 0.5 E\left[Z_{0}^{2}\right]^{1 / 2}=0.5$ and so by Theorem 5.8
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \frac{0.5}{5} \times 5 \times 0.5^{n-1}=0.5^{n}$
So, by iteration $n=7$, the total variation will be less than 0.01 .

In comparison, Fig. 3 shows how the bound compares to a simulated estimate of the total variation distance for this process.

### 5.3 Application to the GARCH $(1,1)$ model

Example $5.12\left(\operatorname{GARCH}(1,1)\right.$ process) Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ be a $\operatorname{GARCH}(1,1)$ process. Then for i.i.d. $Z_{n}$
$X_{n}=\sigma_{n} Z_{n}$
where for $\alpha, \beta, \gamma \in \mathbb{R}$,
$\sigma_{n}^{2}=\alpha^{2}+\beta^{2} X_{n-1}^{2}+\gamma^{2} \sigma_{n-1}^{2}$
See Section 7.3.6 of Doukhan (2018) for more details on this model.

The following theorem provides an upper bound in total variation distance between two copies of the $\operatorname{GARCH}(1,1)$ process.

Theorem 5.13 Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1} \in \mathbb{R}$ be two copies of the GARCH process defined in Example 5.12. Suppose that the density of $Z_{n}$ is centered at 0 and is monotonically decreasing around zero. Then, the process is geometrically ergodic if $\beta^{2} E\left[\left|Z_{0}^{2}\right|\right]+\gamma^{2}<1$. Further suppose that $x_{0}, x_{0}^{\prime}, \sigma_{0}^{2}$, and $\sigma_{0}^{\prime 2}$ are known. Then an upper bound on the total variation distance between the two processes is

$$
\begin{equation*}
\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \frac{D^{n-1}}{\alpha} \sqrt{\beta^{2}\left|x_{0}^{2}-x_{0}^{\prime 2}\right|+\gamma^{2}\left|\sigma_{0}^{2}-\sigma_{0}^{\prime 2}\right|} \tag{18}
\end{equation*}
$$

where $D=\sqrt{\beta^{2} E\left[Z_{0}^{2}\right]+\gamma^{2}}$
Before proving Theorem 5.13, we present some lemmas.
Lemma 5.14 (Contraction condition) The $\operatorname{GARCH}(1,1)$ process satisfies the contraction condition if $D=$ $\sqrt{\beta^{2} E\left[Z_{0}^{2}\right]+\gamma^{2}}<1$ See Sect. D. 3 for a proof.

Lemma 5.15 (Coalescing condition) Suppose that the density of $Z_{n}$ is centred at 0 and is monotonically decreasing

Fig. 2 This figure compares a simulated approximation of $\left\|\mathcal{L}\left(X_{n}^{2}\right)-\mathcal{L}\left(X_{n}^{\prime 2}\right)\right\|$ against the upper bound (Eq. 15). $X_{n}^{2}, X_{n}^{\prime 2}$ are two copies of the LARCH process (i.e.,
$X_{n}^{2}=\left(1+0.5 X_{n-1}^{2}\right) Z_{n}^{2}$ and $\left.Z_{n}^{2} \sim \chi^{2}(1)\right)$ and
$X_{0}^{2}=0.1, X_{0}^{\prime 2}=1.1$. To
simulate total variation, 10 million simulations were run with bin length $=0.01$ for the estimated density function

Fig. 3 This figure compares a simulated approximation of $\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\|$ against the upper bound (Eq. 17). $X_{n}, X_{n}^{\prime}$ are two copies of the asymmetric process (i.e., $X_{n}=$ $\sqrt{\left(0.5 X_{n-1}+3\right)^{2}+5^{2}} Z_{n}, Z_{n} \sim$ $N(0,1)$ ) and $x_{0}=0, x_{0}^{\prime}=5$. To simulate total variation, 10 million simulations were run with bin length $=0.01$ for the estimated density function

Fig. 4 This figure compares a simulated approximation of $\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\|$ against the upper bound (Eq. 19). $X_{n}, X_{n}^{\prime}$ are two copies of the asymmetric process (i.e., $X_{n}=\sigma_{n} Z_{n}$ and $\sigma_{n}^{2}=0.13000+0.1266 X_{n-1}^{2}+$ $0.7922 \sigma_{n-1}^{2}$ and $\left.Z_{n} \sim N(0,1)\right)$ and $X_{0}=0.1, \sigma_{0}=0.01$ and $X_{0}^{\prime}=-0.1, \sigma_{0}^{\prime}=0.1$. To simulate total variation, 1 million simulations were run with bin length $=0.01$ for the estimated density function

Theoretical upper bound vs simulated estimate of total variation distance For the LARCH process


Theoretical upper bound vs simulated estimate of total variation distance
For the asymmetric ARCH process


Theoretical upper bound vs simulated estimate of total variation distance For the GARCH process

around zero. Then the $\operatorname{GARCH}(1,1)$ process satisfies the coalescing condition,
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \frac{D}{\alpha E\left[\left|Z_{0}\right|\right]} E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]$
For $n \geq 2, D=\sqrt{\beta^{2} E\left[Z_{0}^{2}\right]+\gamma^{2}}$. See Sect. D. 3 for a proof.
Lemma 5.16 (Initial condition) Suppose that we know $\sigma_{0}^{2}$, $\sigma_{0}^{\prime 2}$ and $X_{0}, X_{0}^{\prime}$, then
$E\left[\left|X_{1}-X_{1}^{\prime}\right|\right] \leq \sqrt{\beta^{2}\left|X_{0}^{2}-X_{0}^{\prime 2}\right|+\gamma^{2} \mid \sigma_{0}^{2}-\sigma_{0}^{\prime 2}} \mid E\left[\left|Z_{0}\right|\right]$
See Sect. D. 3 for a proof.
Proof of Theorem 5.13 Suppose that the assumptions in Theorem 5.13 are satisfied and let $n \geq 2$. Then the GARCH(1,1) model satisfies the contraction condition (Lemma 5.14) and the coalescing condition (Lemma 5.15). Thus, by the OneShot Coupling Theorem 3.1,
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \frac{D}{\alpha E\left[\left|Z_{0}\right|\right]} D^{n-2} E\left[\left|X_{1}-X_{1}^{\prime}\right|\right]$
Further, by Lemma 5.16 when the initial values $\sigma_{0}^{2}, \sigma_{0}^{\prime 2}, x_{0}$, $x_{0}^{\prime}$ are known,
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \frac{D^{n-1}}{\alpha} \sqrt{\beta^{2}\left|X_{0}^{2}-X_{0}^{\prime 2}\right|+\gamma^{2}\left|\sigma_{0}^{2}-\sigma_{0}^{\prime 2}\right|}$
where $D=\sqrt{\beta^{2} E\left[Z_{0}^{2}\right]+\gamma^{2}}$
Numerical Example 5.17 In Example 10.3.2 of Brockwell and Davis (2002), a GARCH(1,1) model is applied for the daily returns of the Dow Jones Industrial Index between between July 1997 and April 1999. Let

$$
\begin{aligned}
X_{n}= & \sigma_{n} Z_{n} \\
= & \text { excess daily return of the Dow Jones Industrial } \\
& \quad \text { Index at time } n
\end{aligned}
$$

The following is the fitted GARCH volatility estimates when $Z_{n} \sim N(0,1)$,
$\sigma_{n}^{2}=0.13000+0.1266 X_{n-1}^{2}+0.7922 \sigma_{n-1}^{2}$
Suppose that we want to find the total variation of the fitted process with varying initial values representing two market states, $X_{0}=0.1, \sigma_{0}=0.01$ and $X_{0}^{\prime}=-0.1, \sigma_{0}^{\prime}=0.1$ Then by Theorem 5.13,

$$
\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq \sqrt{\frac{0.7922\left|0.01^{2}-0.1^{2}\right|}{0.13}} D^{n-1}
$$

$$
\begin{equation*}
\approx 0.2456 D^{n-1} \tag{19}
\end{equation*}
$$

where $D=\sqrt{0.1266+0.7922}=\sqrt{0.9188}$
By iteration 77, the total variation distance between the two processes is less than 0.01. In comparison, Fig. 4 shows how the bound compares to a simulated estimate of the total variation distance for this process. The actual total variation distance appears to be much smaller than the upper bound.

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## Declaration

Conflict of interest The authors declare that they have no conflict of interest.

## Appendix A: Propositions related to the properties of total variation distance

Proof of Proposition 2.2 Let $\mathcal{A}$ be the sigma field of $\mathcal{X}$ and $\mathcal{B}$ be the sigma field of $\mathcal{Y}$.

First note that $f^{-1}(\mathcal{B})=\left\{f^{-1}(B): B \in \mathcal{B}\right\}=\mathcal{A}$ :

- $f^{-1}(\mathcal{B}) \subset \mathcal{A}$ : For $B \in \mathcal{B}, f^{-1}(B) \subset \mathcal{A}$ by measurability.
- $\mathcal{A} \subset f^{-1}(\mathcal{B})$ : Let $A \in \mathcal{A}$. Then $f(A) \in \mathcal{B}$ and $f^{-1}(f(A)) \in f^{-1}(\mathcal{B})$ by definition. By invertibility, $f^{-1}(f(A))=A$ and so $A \in f^{-1}(\mathcal{B})$.

The equality in Eq. 2 can then be proven as follows,

$$
\begin{aligned}
& \left\|\mathcal{L}(f(X))-\mathcal{L}\left(f\left(X^{\prime}\right)\right)\right\| \\
& \quad=\sup _{B \in f(\mathcal{B})}\left|P(f(X) \in B)-P\left(f\left(X^{\prime}\right) \in B\right)\right| \\
& \quad=\sup _{B \in f(\mathcal{B})}\left|P\left(X \in f^{-1}(B)\right)-P\left(X^{\prime} \in f^{-1}(B)\right)\right| \\
& \quad=\sup _{A \in \mathcal{A}}\left|P(X \in A)-P\left(X^{\prime} \in A\right)\right| \text { Since } f^{-1}(\mathcal{B})=\mathcal{A} \\
& \quad=\left\|\mathcal{L}(X)-\mathcal{L}\left(X^{\prime}\right)\right\|
\end{aligned}
$$

## Proof of Proposition 2.3

$$
\begin{aligned}
& \left\|\mathcal{L}(X)-\mathcal{L}\left(X^{\prime}\right)\right\|=\sup _{A \in \mathcal{B}}\left|P(X \in A)-P\left(X^{\prime} \in A\right)\right| \\
& \quad=\sup _{A \in \mathcal{B}} \int_{\mathcal{Y}} P(X \in A \mid y)-P\left(X^{\prime} \in A \mid y\right) \mu(d y) \mid \\
& \quad \leq \sup _{A \in \mathcal{B}} \int_{\mathcal{Y}}\left|P(X \in A \mid y)-P\left(X^{\prime} \in A \mid y\right)\right| \mu(d y)
\end{aligned}
$$

by Jensen's inequality

$$
\begin{aligned}
& \leq \int_{\mathcal{Y}} \sup _{A \in \mathcal{B}}\left|P(X \in A \mid y)-P\left(X^{\prime} \in A \mid y\right)\right| \mu(d y) \\
& \leq E\left[\left\|\mathcal{L}(X \mid Y)-\mathcal{L}\left(X^{\prime} \mid Y\right)\right\|\right]
\end{aligned}
$$

Proof of Proposition 2.4 To prove this, we use the concept of maximal coupling over the coordinates. By maximal coupling, for $i \in\{1, \ldots, d\}$ there exists random variables $X_{i, n}^{M}, X_{i, n}^{\prime M}$ such that $X_{i, n} \stackrel{d}{=} X_{i, n}^{M}$ and $X_{i, n}^{\prime} \stackrel{d}{=} X_{i, n}^{\prime M}$ and
$\left\|\mathcal{L}\left(X_{i, n}\right)-\mathcal{L}\left(X_{i, n}^{\prime}\right)\right\|=P\left(X_{i, n}^{M} \neq X_{i, n}^{\prime}{ }^{M}\right)$
(see Proposition 3g of Roberts and Rosenthal (2004) or Section 2 of Böttcher (2017)).

Further, there exists a unique product measure such that for any $A_{1}, \ldots A_{d} \in \mathcal{B}, P\left(\cap_{i=1}^{d}\left[X_{i, n}^{M} \in A_{i}\right]\right)=\prod_{i=1}^{d} P\left(X_{i, n}^{M} \in\right.$ $A_{i}$ ) (Theorem 18.2 of Billingsley (2012)). For the unique product measure, the following equality holds,

$$
\begin{aligned}
& P\left(\cap_{i=1}^{d} X_{i, n}^{M} \in A_{i}\right)=\prod_{i=1}^{d} P\left(X_{i, n}^{M} \in A_{i}\right) \\
& \quad=\prod_{i=1}^{d} P\left(X_{i, n} \in A_{i}\right)=P\left(\cap_{i=1}^{d} X_{i, n} \in A_{i}\right)
\end{aligned}
$$

And so by uniqueness, for $A \in \mathcal{B}^{\mathrm{d}}, P\left(X_{n}^{M} \in A\right)=P\left(X_{n} \in\right.$ A). By definition, this means that $\vec{X}_{n} \stackrel{d}{=} \vec{X}_{n}^{M}$, which implies that $\left(\vec{X}_{n}^{M}, \vec{X}_{n}^{\prime M}\right) \in \mathcal{C}\left(\vec{X}_{n}, \vec{X}_{n}^{\prime}\right)$, the set of all couplings of $\vec{X}_{n}, \vec{X}_{n}^{\prime}$.

We now use $\vec{X}_{n}^{M}, \vec{X}_{n}^{\prime M}$ to prove Eq. 3 .

$$
\begin{aligned}
& \left\|\mathcal{L}\left(\vec{X}_{n}\right)-\mathcal{L}\left(\vec{X}_{n}^{\prime}\right)\right\| \\
& \quad=\inf _{\vec{Y}, \vec{Y}^{\prime} \in \mathcal{C}\left(\vec{X}_{n}, \vec{X}_{n}^{\prime}\right)} P\left(\vec{Y} \neq \vec{Y}^{\prime}\right) \text { by Eq. } 2.4 \text { of Böttcher }(2017) \\
& \quad \leq P\left(\vec{X}_{n}^{M} \neq \vec{X}_{n}^{\prime M}\right) \\
& \quad=P\left(\cup_{i=1}^{d}\left[X_{i, n}^{M} \neq X_{i, n}^{\prime M}\right]\right) \\
& \quad \leq \sum_{i=1}^{d} P\left(X_{i, n}^{M} \neq X_{i, n}^{\prime M}\right) \quad \text { by subadditivity } \\
& \quad \leq d A r^{n}
\end{aligned}
$$

## Appendix B: Lemmas related to the Sideways Theorem

The following are lemmas and corresponding proofs and corollaries related to the Sideways Theorem (4.2).

## B.0.1 Lemmas providing an upper bound on the integral difference between a function and a corresponding shift

The following lemmas are used in the proof of Lemma 4.3.
Lemma B. 1 For any invertible, continuousfunction $f: \mathbb{R} \rightarrow$ $\mathbb{R}$ where the codomain is $f(\mathbb{R})=(a, b)$ and $\Delta>0$,

$$
\int_{\mathbb{R}}|f(x+\Delta)-f(x)| d x=(b-a) \Delta
$$

Proof Since $f$ is invertible and continuous, it is strictly monotone (Lemma 3.8 if Hairer and Wanner (2008)). Assume that $f$ is strictly increasing. The integral can be written as follows,

$$
\begin{aligned}
\int_{\mathbb{R}} & |f(x+\Delta)-f(x)| d x=\int_{\mathbb{R}} f(x+\Delta)-f(x) d x \\
& =\int_{\mathbb{R}} \int_{a}^{b} I_{f(x+\Delta)<y<f(x)} d y d x \\
& =\int_{\mathbb{R}} \int_{a}^{b} I_{f^{-1}(y)-\Delta<x<f^{-1}(y)} d y d x \\
& =\int_{a}^{b} \int_{\mathbb{R}} I_{f^{-1}(y)-\Delta<x<f^{-1}(y)} d x d y \text { by Fubini’s Theorem } \\
& =\int_{a}^{b} \Delta d y \\
& =(b-a) \Delta
\end{aligned}
$$

If $f$ is strictly decreasing apply the transform $h(x)=$ $a+b-f(x)$. The function $h$ is a strictly increasing invertible function with codomain $(a, b)$ and so using the previous result for increasing functions,

$$
\begin{aligned}
& \int_{\mathbb{R}}|f(x+\Delta)-f(x)| d x \\
& \quad=\int_{\mathbb{R}}|h(x+\Delta)-h(x)| d x=(b-a) \Delta
\end{aligned}
$$

Lemma B. 2 Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be a continuous function that is invertible over the set $(c, d)$ and is a constant function over $(c, d)^{C}$. Further suppose that the codomain is $f(\mathbb{R})=$ $(a, b)$. Then for $\Delta>0$, we get that
$\int_{\mathbb{R}}|f(x+\Delta)-f(x)| d x=(b-a) \Delta$

Proof Assume that $f$ is an increasing function and so $f(c)=$ $a, f(d)=b$ and $|f(x+\Delta)-f(x)|=f(x+\Delta)-f(x)$.

Let $0<\epsilon<(c-d) / 2$ and define
$g_{\epsilon}(x)= \begin{cases}(f(c+\epsilon)-a)\left(1-e^{x-c-\epsilon}\right)+a & \text { whenx } \in(-\infty, c+\epsilon] \\ f(x) & \text { whenx } \in(c+\epsilon, d-\epsilon] \\ (f(d-\epsilon)-b)\left(1-e^{d-\epsilon-x}\right)+b \text { whenx } \in(d-\epsilon, \infty)\end{cases}$
Note that $g_{\epsilon}(x)$ is continuous, invertible, an increasing function and the codomain is $(a, b)$. By Lemma B. 1 for each $\epsilon>0$
$\int_{\mathbb{R}} g_{\epsilon}(x+\Delta)-g_{\epsilon}(x) d x=(b-a) \Delta$
Further, for all $x \in \mathbb{R}, \lim _{\epsilon \rightarrow 0} g_{\epsilon}(x+\Delta)-g_{\epsilon}(x)=$ $f(x+\Delta)-f(x)$ and so $g_{\epsilon}(x+\Delta)-g_{\epsilon}(x)$ converges pointwise to $f(x+\Delta)-f(x)$. Next, for $0<\epsilon<(c-d) / 2$, $\left|g_{\epsilon}(x+\Delta)-g_{\epsilon}(x)\right|<2|b|$ and so the function $g_{\epsilon}(x+\Delta)-$ $g_{\epsilon}(x)$ is uniformly bounded. The above statements allow us to apply the dominated convergence Theorem (Theorem 16.5 of Billingsley (2012)) and so

$$
\begin{aligned}
& \int_{\mathbb{R}} f(x+\Delta)-f(x) d x \\
& \quad=\lim _{\epsilon \rightarrow 0} \int_{\mathbb{R}} g_{\epsilon}(x+\Delta)-g_{\epsilon}(x) d x=(b-a) \Delta
\end{aligned}
$$

If $f$ is strictly decreasing apply the transform $h(x)=$ $a+b-f(x)$. The function $h$ is a strictly increasing invertible function with codomain $(a, b)$ and so using the previous result for increasing functions,

$$
\begin{aligned}
& \int_{\mathbb{R}}|f(x+\Delta)-f(x)| d x=\int_{\mathbb{R}} \mid h(x+\Delta) \\
& \quad-h(x) \mid d x=(b-a) \Delta
\end{aligned}
$$

Lemma B. 3 Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be a continuous function with the following properties:

- the codomain is $(0, K)$
- $\left(m_{1}, m_{2}, \ldots, m_{M}\right)$ are the local maxima and minima points
- $\lim _{x \rightarrow \infty} f(x)=0$ and $\lim _{x \rightarrow-\infty} f(x)=0$

Further suppose that $\Delta<\max _{i=2, \ldots, M}\left\{m_{i}-m_{i-1}\right\}$. Then
$\int_{\mathbb{R}}|f(x-\Delta)-f(x)| d x \leq K(M+1) \Delta$
Proof Since $\Delta<\max _{i=2, \ldots, M}\left\{m_{i}-m_{i-1}\right\}$, we have that $m_{1}-\Delta<m_{1}<m_{2}-\Delta<\ldots<m_{M}$. Let $I_{1}, \ldots, I_{M}$ be the intersection points or the points where $f\left(I_{i}\right)=f\left(I_{i}-\Delta\right)$.

Show that $m_{i}-\Delta<I_{i}<m_{i}$ : Suppose that $m_{i}$ is a local maximum point. Let $g(x)=f(x+\Delta)$. Within the interval
$\left(m_{i}-\Delta, m_{i}\right), f^{\prime}(x)>0$ and $g^{\prime}(x)<0$ by assumption. This implies that $f\left(m_{i}-\Delta\right)<f\left(m_{i}\right)$ and $g\left(m_{i}-\Delta\right)>g\left(m_{i}\right)$ by the Mean Value Theorem. Further since $g\left(m_{i}-\Delta\right)=f\left(m_{i}\right)$ we have that $g\left(m_{i}-\Delta\right)>f\left(m_{i}-\Delta\right)$ and $g\left(m_{i}\right)<f\left(m_{i}\right)$.

Let $h(x)=g(x)-f(x)$. Then $h\left(m_{i}-\Delta\right)>0$ and $h\left(m_{i}\right)<0$ further $h$ is a strictly decreasing function over ( $m_{i}-\Delta, m_{i}$ ) since $g,-f$ are strictly decreasing functions over the same interval. So by the intermediate value theorem, there exists an $\xi \in\left(m_{i}-\Delta, m_{i}\right)$ such that $h(\xi)=0$ or $f(\xi)=g(\xi)=f(\xi+\Delta)$. Further by injectivity, $\xi$ is unique. Let $I_{i}=\xi$. A similar proof can be given for when $m_{i}$ is a local minimum.

Show that $\int_{I_{i}}^{I_{i+1}}|f(x+\Delta)-f(x)| d x \leq K \Delta$ : Note first that $m_{i}-\Delta<I_{i}<m_{i}<m_{i+1}-\Delta<I_{i+1}<m_{i+1}$ further define
$f_{i}(x)= \begin{cases}f\left(m_{i}\right) & \text { when } x \in\left(-\infty, m_{i}\right] \\ f(x) & \text { when } x \in\left(m_{i}, m_{i+1}\right] \\ f\left(m_{i+1}\right) & \text { when } x \in\left(m_{i+1}, \infty\right)\end{cases}$
Note that over the interval ( $m_{i}, m_{i+1}$ ], the function $f$ is either a strictly increasing or a strictly decreasing function.

$$
\begin{aligned}
& \int_{I_{i}}^{I_{i+1}}|f(x+\Delta)-f(x)| d x \\
&= \int_{I_{i}}^{m_{i}}|f(x+\Delta)-f(x)| d x \\
&+\int_{m_{i}}^{m_{i+1}-\Delta}|f(x+\Delta)-f(x)| d x \\
&+\int_{m_{i+1}-\Delta}^{I_{i+1}}|f(x+\Delta)-f(x)| d x \\
& \leq \int_{I_{i}}^{m_{i}}\left|f(x+\Delta)-f\left(m_{i}\right)\right| d x \\
&+\int_{m_{i}}^{m_{i+1}-\Delta}|f(x+\Delta)-f(x)| d x \\
&+\int_{m_{i+1}-\Delta}^{I_{i+1}}\left|f\left(m_{i+1}\right)-f(x)\right| d x \\
&= \int_{I_{i}}^{m_{i}}\left|f_{i}(x+\Delta)-f_{i}(x)\right| d x \\
&+\int_{m_{i}}^{m_{i+1}-\Delta}\left|f_{i}(x+\Delta)-f_{i}(x)\right| d x \\
&+\int_{m_{i+1}-\Delta}^{I_{i+1}}\left|f_{i}(x+\Delta)-f_{i}(x)\right| d x \\
&= \int_{I_{i}}^{I_{i+1}}\left|f_{i}(x+\Delta)-f_{i}(x)\right| d x \\
& \leq \int_{m_{i}-\Delta}^{m_{i+1}}\left|f_{i}(x+\Delta)-f_{i}(x)\right| d x \\
&= \int_{\mathbb{R}}\left|f_{i}(x+\Delta)-f_{i}(x)\right| d x \\
&
\end{aligned}
$$

$$
=\left|f\left(m_{i}\right)-f\left(m_{i+1}\right)\right| \Delta \leq K \Delta
$$

The last equality is a result of Lemma B.2.
By similar reasoning, it can be shown that

$$
\begin{array}{r}
\int_{-\infty}^{I_{1}}|f(x+\Delta)-f(x)| d x \leq K \Delta \\
\int_{I_{M}}^{\infty}|f(x+\Delta)-f(x)| d x \leq K \Delta
\end{array}
$$

Finally note that the intersection points partition $\mathbb{R}$ into $M+1$ subsets and so

$$
\int_{\mathbb{R}}|f(x-\Delta)-f(x)| d x \leq K(M+1) \Delta
$$

## Proof of Lemma 4.3

Lemma 4.3 represents the coalescing condition for the Sideways Theorem 4.2.

Proof of Lemma 4.3 Set $\theta_{1, n}=\theta_{1, n}^{\prime}$. Define
$\Delta=g\left(\theta_{1, n}, X_{n-1}\right)-g\left(\theta_{1, n}, X_{n-1}^{\prime}\right)$
Let $f_{X_{n}}, f_{X_{n}^{\prime}}$ be the density functions for $X_{n}, X_{n}^{\prime}$, respectively, and $f_{\theta_{2, n}}, f_{\theta_{2, n}+\Delta}$ be the density functions for $\theta_{2, n}, \theta_{2, n}$ $+\Delta$.

Suppose that $\Delta, X_{n-1}, X_{n-1}^{\prime} \in \mathbb{R}$ are known and so,

$$
\begin{aligned}
X_{n} & =g\left(\theta_{1, n}, X_{n-1}\right)+\theta_{2, n} \Longrightarrow \theta_{2, n} \\
& =X_{n}-g\left(\theta_{1, n}, X_{n-1}\right) \\
X_{n}^{\prime} & =g\left(\theta_{1, n}, X_{n-1}^{\prime}\right)+\theta_{2, n}^{\prime} \Longrightarrow \theta_{2, n}^{\prime}-\Delta \\
& =X_{n}^{\prime}-g\left(\theta_{1, n}, X_{n-1}\right)
\end{aligned}
$$

We know that $\theta_{2, n} \stackrel{d}{=} \theta_{2, n}^{\prime}$ and in general $\Delta, \theta_{1, n}$ are random variables, so

$$
\begin{aligned}
& \left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \\
& \quad \leq E_{\theta_{1, n}, \Delta}\left[\left\|\mathcal{L}\left(X_{n} \mid \theta_{1, n}, \Delta\right)-\mathcal{L}\left(X_{n}^{\prime} \mid \theta_{1, n}, \Delta\right)\right\|\right]
\end{aligned}
$$

by Proposition2.3
$=E_{\theta_{1, n}, \Delta}\left[\left\|\mathcal{L}\left(\theta_{2, n} \mid \theta_{1, n}\right)-\mathcal{L}\left(\theta_{2, n}-\Delta \mid \theta_{1, n}\right)\right\|\right]$ by Proposition 2.3

By the assumptions in the theorem, the density of $\theta_{2, n}$ is continuous with $M$ extrema points and has a codomain that is in $(0, K)$. Let $\left(m_{1}, m_{2}, \ldots, m_{M}\right)$ be the local extrema points where $m_{i}<m_{j}$ if $i<j$ and $L \leq \max _{2 \leq i \leq M}\left\{m_{i}-m_{i-1}\right\}$ be the maximum distance between two local extrema points. So, continuing from the inequality B1 and by the definition of total variation, Eq. 1,
$\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\|$

$$
\begin{aligned}
& \leq E_{\theta_{1, n}}\left[E_{\Delta}\left[\frac{1}{2} \int_{\mathbb{R}}\left|f_{\theta_{2, n}}\left(x \mid \theta_{1, n}\right)-f_{\theta_{2, n}-\Delta}\left(x \mid \theta_{1, n}\right)\right| d x\right]\right] \\
& =E_{\theta_{1, n}}\left[E_{\Delta}\left[\frac{1}{2} \int_{\mathbb{R}}\left|f_{\theta_{2, n}}\left(x \mid \theta_{1, n}\right)-f_{\theta_{2, n}}\left(x+\Delta \mid \theta_{1, n}\right)\right| d x\right]\right] \\
& =E_{\theta_{1, n}}\left[E_{\Delta}\left[\frac{1}{2} \int_{\mathbb{R}}\left|f_{\theta_{2, n}}\left(x \mid \theta_{1, n}\right)-f_{\theta_{2, n}}\left(x+\Delta \mid \theta_{1, n}\right)\right| d x I_{\Delta<L}\right]\right] \\
& +E_{\theta_{1, n}}\left[E_{\Delta}\left[\frac{1}{2} \int_{\mathbb{R}}\left|f_{\theta_{2, n}}\left(x \mid \theta_{1, n}\right)-f_{\theta_{2, n}}\left(x+\Delta \mid \theta_{1, n}\right)\right| d x I_{\Delta>L}\right]\right] \\
& \leq \frac{1}{2} E_{\theta_{1, n}}\left[E_{\Delta}[K(M+1)|\Delta|]\right]+P_{\Delta}(|\Delta|>L) \\
& \leq \frac{K(M+1)}{2} E_{\Delta}[|\Delta|]+\frac{E_{\Delta}[|\Delta|]}{L}
\end{aligned}
$$

The second last inequality is a result of Lemma B.3. The coalescing condition is thus satisfied as follows with $C=$ $\frac{K(M+1)}{2}+\frac{I_{M>1}}{L}$,

$$
\begin{aligned}
& \left\|\mathcal{L}\left(X_{n+1}\right)-\mathcal{L}\left(X_{n+1}^{\prime}\right)\right\| \\
& \quad \leq C E\left[\left|g\left(\theta_{1, n}, X_{n-1}\right)-g\left(\theta_{1, n}, X_{n-1}^{\prime}\right)\right|\right] \\
& \quad=C E\left[\mid g\left(\theta_{1, n}, X_{n-1}\right)+\theta_{2, n}\right. \\
& \left.\quad-\left(g\left(\theta_{1, n}, X_{n-1}^{\prime}\right)+\theta_{2, n}\right) \mid\right] \\
& \quad=C E\left[\left|X_{n}-X_{n}^{\prime}\right|\right]
\end{aligned}
$$

## Appendix C: Lemmas for random-functional autoregressive process examples

## C. 1 Proof of Lemma 4.5

Proof of Lemma 4.5 First note that

$$
\begin{aligned}
& E\left[\left|X_{n+2}-X_{n+2}^{\prime}\right| \mid X_{n}=x, X_{n}^{\prime}=y\right] \\
& \quad=E\left[\left|g\left(\frac{1}{2}(x-\sin x)+Z_{n}\right)-g\left(\frac{1}{2}(y-\sin y)+Z_{n}\right)\right|\right] \\
& \quad=\frac{1}{2} E\left[\left\lvert\, \frac{1}{2}(x-y+\sin y-\sin x)\right.\right. \\
& \left.\left.\quad+\sin \left(\frac{1}{2}(y-\sin y)+Z_{n}\right)-\sin \left(\frac{1}{2}(x-\sin x)+Z_{n}\right) \right\rvert\,\right] \\
& \quad=\frac{1}{2} E[|g(x, y)+G(x, y)|]
\end{aligned}
$$

where $g(x, y)=\frac{1}{2}(x-y+\sin y-\sin x)$ and $G(x, y)=$ $\sin \left(\frac{1}{2}(y-\sin y)+Z_{n}\right)-\sin \left(\frac{1}{2}(x-\sin x)+Z_{n}\right)$. By trigonometric identities ${ }^{1}$, for $k(x, y)=\frac{x+y-\sin y-\sin x}{4}$ and $h(x, y)=\frac{y-x+\sin x-\sin y}{4}$.
$G(x, y)=2 \cos \left(\frac{x+y-\sin y-\sin x}{4}+Z_{n}\right)$

[^1]\[

$$
\begin{aligned}
& \sin \left(\frac{y-x+\sin x-\sin y}{4}\right) \\
= & 2 \cos \left(k(x, y)+Z_{n}\right) \sin h(x, y) \\
= & 2 \sin h(x, y)\left(\cos Z_{n} \cos k(x, y)\right. \\
& \left.+\sin Z_{n} \sin k(x, y)\right)
\end{aligned}
$$
\]

Let $X_{n}=\frac{Z_{n}^{2}}{C}, Y_{n}=\frac{C}{2 G_{n}}$. We can rewrite $\sigma_{n}^{2}=$ $X_{n} Y_{n} \sigma_{n-1}^{2}+Y_{n}$ where $X_{n} \sim \Gamma\left(\frac{p}{2}, \frac{C}{2}\right)$ and $Y_{n} \sim \Gamma^{-1}$ $\left(\frac{k+p}{2}, \frac{C}{2}\right)$. Using the notation from the Sideways Theorem $4.2 \theta_{1, n}=X_{n} Y_{n}$ and $\theta_{2, n}=Y_{n}$.

Since $\beta_{n}$ can be written as a random function of $\sigma_{n}^{2}$,

And so,

$$
\begin{aligned}
& E\left[\left|X_{n+2}-X_{n+2}^{\prime}\right| \mid X_{n}=x, X_{n}^{\prime}=y\right] \\
& =\frac{1}{2} E\left[\left|g(x, y)+2 \sin h(x, y)\left(\cos Z_{n} \cos k(x, y)+\sin Z_{n} \sin k(x, y)\right)\right|\right] \\
& \leq \frac{1}{2} \sqrt{E\left[\left(g(x, y)+2 \sin h(x, y)\left(\cos Z_{n} \cos k(x, y)+\sin Z_{n} \sin k(x, y)\right)\right)^{2}\right]} \\
& =\frac{1}{2} \sqrt{g(x, y)^{2}+4 \frac{g(x, y) \sin h(x, y) \cos k(x, y)}{e^{1 / 2}}+2 \sin ^{2} h(x, y)\left(1+\frac{\cos ^{2} k(x, y)-\sin ^{2} k(x, y)}{e^{2}}\right)} \\
& =\frac{1}{\sqrt{2}} \sqrt{2 h(x, y)^{2}-4 \frac{h(x, y) \sin h(x, y) \cos k(x, y)}{e^{1 / 2}}+\sin ^{2} h(x, y)\left(1+\frac{\cos ^{2} k(x, y)-\sin ^{2} k(x, y)}{e^{2}}\right)}
\end{aligned}
$$

## C. 2 Proof of lemmas used in Theorem 4.8

To prove the first part of this theorem, we apply the deinitialization technique which shows how the convergence rate of a Markov chain can be bounded above by the convergence rate of a more simpler Markov chain that includes sufficient information on the Markov chain of interest. The concept of de-initialization and a proposition that bounds total variation is provided below.

Definition C. 1 (De-initialization) Let $\left\{X_{n}\right\}_{n \geq 1}$ be a Markov chain. A Markov chain $\left\{Y_{n}\right\}_{n \geq 1}$ is a de-initialization of $\left\{X_{n}\right\}_{n \geq 1}$ if for each $n \geq 1$
$\mathcal{L}\left(X_{n} \mid X_{0}, Y_{n}\right)=\mathcal{L}\left(X_{n} \mid Y_{n}\right)$
Proposition C. 2 (Theorem 1 of Roberts and Rosenthal (2001)) Let $\left\{Y_{n}\right\}_{n \geq 1}$ be a de-initialization of $\left\{X_{n}\right\}_{n \geq 1}$ thenfor any two initial distributions $X_{0} \sim \mu$ and $X_{0}^{\prime} \sim \mu^{\prime}$,

$$
\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq\left\|\mathcal{L}\left(Y_{n}\right)-\mathcal{L}\left(Y_{n}^{\prime}\right)\right\|
$$

Proof of Lemma 4.9 Note that $\beta_{n}=\tilde{\beta}+\sigma_{n-1} Z_{n}, Z_{n} \sim$ $N_{p}\left(0, A^{-1}\right)$ can be written as a random function of $\sigma_{n}^{2}$. Substituting $\beta_{n}, \sigma_{n}^{2}$ can then be written as a random function of its previous value for independent $Z_{n}^{2} \sim \chi^{2}(p)$ and $G_{n} \sim \Gamma\left(\frac{k+p}{2}, 1\right)$,
$\sigma_{n}^{2}=\frac{Z_{n}^{2}}{C} \frac{C}{2 G_{n}} \sigma_{n-1}^{2}+\frac{C}{2 G_{n}}$
$\mathcal{L}\left(\beta_{n}, \sigma_{n}^{2} \mid \beta_{0}, \sigma_{0}^{2}, \sigma_{n}^{2}\right)=\mathcal{L}\left(\beta_{n}, \sigma_{n}^{2} \mid \sigma_{n}^{2}\right)$
and so $\sigma_{n}^{2}$ is a de-initialization of $\left(\beta_{n}, \sigma_{n}^{2}\right)$. By Proposition C.2,

$$
\left\|\mathcal{L}\left(\beta_{n}, \sigma_{n}^{2}\right)-\mathcal{L}\left(\beta_{n}^{\prime}, \sigma_{n}^{\prime 2}\right)\right\| \leq\left\|\mathcal{L}\left(\sigma_{n}^{2}\right)-\mathcal{L}\left(\sigma_{n}^{\prime 2}\right)\right\|
$$

We are thus interested in evaluating the convergence rate of $\sigma_{n}^{2}$ to bound the convergence rate of $\left(\beta_{n}, \sigma_{n}^{2}\right)$.

To interpret this in another way, if $\sigma_{n}^{2}$ couples then the distribution of $\beta_{n}$ is the same for both iterations, so it is automatically coupled. An alternative proof can be made using the results from Liu et al. (1994).

Proof of Lemma 4.10 By Lemma 4.9, $\theta_{1, n}=X_{n} Y_{n}$ and so,

$$
\begin{aligned}
K= & E\left[\left|\theta_{1, n}\right|\right]=E\left[X_{n} Y_{n}\right]=E\left[X_{n}\right] E\left[Y_{n}\right] \\
& =\frac{p}{C} \frac{C}{k+p-2}=\frac{p}{k+p-2}
\end{aligned}
$$

Proof of Lemma 4.11 Calculate the conditional density $\theta_{2, n} \mid \theta_{1, n}$ We remove the subscript $n$ on the random variables. Let $X, Y$ be as described in Lemma 4.9. Since the random variables are independent, the joint density is the product of the densities.

$$
\begin{align*}
f_{X, Y}(x, y)= & \frac{C / 2}{\Gamma(p / 2)} x^{p / 2-1} e^{x C / 2} \\
& \frac{C / 2}{\Gamma((k+p) / 2)} y^{-(k+p) / 2-1} e^{-\frac{C / 2}{y}} \tag{C3}
\end{align*}
$$

Then $\left(\theta_{1}, \theta_{2}\right)=(X Y, Y)$ is a transformation with the Jacobian $|J|=\theta_{2}^{-1}$ and the density written as follows,

$$
\begin{aligned}
f_{\theta_{1}, \theta_{2}}\left(\theta_{1}, \theta_{2}\right)= & f_{X, Y}\left(\frac{\theta_{1}}{\theta_{2}}, \theta_{2}\right) \theta_{2}^{-1} \\
= & \frac{C / 2}{\Gamma(p / 2)}\left(\frac{\theta_{1}}{\theta_{2}}\right)^{p / 2-1} e^{-\frac{\theta_{1}}{\theta_{2}} C / 2} \\
& \frac{C / 2}{\Gamma((k+p) / 2)} \theta_{2}^{-(k+p) / 2-1} e^{-\frac{C / 2}{\theta_{2}} \theta_{2}^{-1}}
\end{aligned}
$$

Next $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ is proportional to $f_{\theta_{1}, \theta_{2}}\left(\theta_{1}, \theta_{2}\right)$ and so we can derive the conditional density of $\theta_{2}$ as follows,

$$
\begin{align*}
f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right) & \propto f_{\theta_{1}, \theta_{2}}\left(\theta_{1}, \theta_{2}\right)  \tag{C4}\\
& \propto \theta_{2}^{1-p / 2} e^{-\frac{1}{\theta_{2}} \theta_{1} C / 2} \theta_{2}^{-(k+p) / 2-1} e^{-\frac{1}{\theta_{2}} C / 2} \theta_{2}^{-1} \\
& =\theta_{2}^{-(p / 2+(k+p) / 2)-1} e^{-\frac{1}{\theta_{2}}\left(\theta_{1}+1\right) C / 2} \tag{C5}
\end{align*}
$$

This is proportional to an inverse gamma distribution and so, $\theta_{2} \left\lvert\, \theta_{1} \sim \Gamma^{-1}\left(\frac{k+2 p}{2},\left(\theta_{1}+1\right) C / 2\right)\right.$. Since the conditional density is an inverse gamma distribution, the number of modes is $M=1$ and the density function is continuous.

Calculate the maximum value of $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ : Fig. 5 shows how the maximum value of the density increases as the shape, $\left(\theta_{1}+1\right) C / 2$ decreases when the rate, $\frac{k+2 p}{2}$ is fixed. It can also be shown from equation C 4 that the density function of $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ is maximized when $\theta_{1}=0$ since the normalizing constant will be the largest. This means that $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ reaches its maximum height when $\theta_{1}=0$ and so we find the value of $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ evaluated at $\theta_{2}=$ $\frac{C}{k+2 p+2}$, the mode (Section 5.3 of Hoff (2009)).

$$
\begin{aligned}
K & =f_{\theta_{2} \mid \theta_{1}}\left(\left.\frac{C}{k+2 p+2} \right\rvert\, \theta_{1}=0\right) \\
& =\left.\frac{(C / 2)^{\frac{k+2 p}{2}}}{\Gamma\left(\frac{k+2 p}{2}\right)} y^{-\frac{k+2 p}{2}-1} e^{-\frac{C / 2}{y}}\right|_{y=\frac{C}{k+2 p+2}} \\
& =\frac{(C / 2)^{\frac{k+2 p}{2}}}{\Gamma\left(\frac{k+2 p}{2}\right)}\left(\frac{C}{k+2 p+2}\right)^{-\frac{k+2 p}{2}-1} e^{-\frac{k+2 p+2}{2}} \\
& =\frac{(C / 2)^{\frac{k+2 p}{2}}}{\Gamma\left(\frac{k+2 p}{2}\right)}\left(\frac{k+2 p+2}{C}\right)^{\frac{k+2 p}{2}+1} e^{-\frac{k+2 p+2}{2}}
\end{aligned}
$$

And so,
$K=\frac{(C / 2)^{\frac{k+2 p}{2}}}{\Gamma\left(\frac{k+2 p}{2}\right)}\left(\frac{k+2 p+2}{C}\right)^{\frac{k+2 p}{2}+1} e^{-\frac{k+2 p+2}{2}}$

## C. 3 Proof of lemmas used in Theorem 4.14

Proof of Lemma 4.15 The iteration $\tau_{n+1}^{-1}$ can be written as a function of its previous value, $\tau_{n}^{-1}$ since $\mu_{n+1}=\bar{y}+$ $Z_{n+1} / \sqrt{J \tau_{n}}$.
$\tau_{n+1}^{-1}=\frac{Z_{n+1}^{2}}{S} \frac{S}{2 G_{n+1}} \tau_{n}^{-1}+\frac{S}{2 G_{n+1}}$
Next we can rewrite, $\tau_{n}^{-1}=X_{n} Y_{n} \tau_{n-1}^{-1}+Y_{n}$ where $X_{n}=$ $\frac{Z_{t+1}^{2}}{S} \sim \Gamma\left(\frac{1}{2}, \frac{S}{2}\right)$ and $Y_{n}=\frac{S}{2 G_{t+1}} \sim \Gamma^{-1}\left(\frac{J+2}{2}, \frac{S}{2}\right)$.

Since $\left(\mu_{n}, \tau_{n}^{-1}\right)$ can be written as a random function of $\tau_{n}^{-1}$,
$\mathcal{L}\left(\mu_{n}, \tau_{n}^{-1} \mid \mu_{0}, \tau_{0}^{-1}, \tau_{n}^{-1}\right)=\mathcal{L}\left(\mu_{n}, \tau_{n}^{-1} \mid \tau_{n}^{-1}\right)$
and $\tau_{n}^{-1}$ is a de-initialization of $\left(\mu_{n}, \tau_{n}^{-1}\right)$. Further, by Proposition C.2,
$\left\|\mathcal{L}\left(\mu_{n}, \tau_{n}^{-1}\right)-\mathcal{L}\left(\mu_{n}^{\prime}, \tau_{n}^{\prime-1}\right)\right\| \leq\left\|\mathcal{L}\left(\tau_{n}^{-1}\right)-\mathcal{L}\left(\tau_{n}^{\prime-1}\right)\right\|$
To interpret this in another way, if $\tau_{n}$ couples then the distribution of $\mu_{n}$ is the same for both iterations, so it is automatically coupled. An alternative proof can be made using the results from Liu et al. (1994).

Proof of Lemma 4.16 By Lemma 4.15, $\theta_{1, n}=X_{n} Y_{n}$ and so by Corollary 4.6
$D=E\left[\left|\theta_{1, n}\right|\right]=E\left[X_{n} Y_{n}\right]=E\left[X_{n}\right] E\left[Y_{n}\right]=\frac{1}{S} \frac{S}{J}=\frac{1}{J}$

Proof of Lemma 4.17 To find $M, K$ and show that the conditional density is continuous, we (a) show that $\theta_{2} \mid \theta_{1} \sim$ $\Gamma^{-1}\left(\frac{J-1}{2},\left(\theta_{1}+1\right) S / 2\right)$, which directly implies that the conditional distribution is continuous and $M=1$ and we (b) we find the value of $K$.
(a) Calculate the conditional density $\theta_{2, n} \mid \theta_{1, n}$ For simplicity, we remove the subscript $n$ on the random variables. Let $X, Y$ be as described in Lemma 4.15. Since the random variables are independent, the joint density is the product of the densities.

$$
\begin{align*}
& f_{X, Y}(x, y)= \frac{S / 2}{\Gamma(1 / 2)} x^{1 / 2-1} e^{x S / 2} \frac{S / 2}{\Gamma((J+2) / 2)} \\
& y^{-(J+2) / 2-1} e^{-\frac{S / 2}{y}} \tag{C9}
\end{align*}
$$

Fig. 5 Inverse gamma density
when $\alpha=100$ and $\beta=1,10,100$

## Inverse gamma density for different rates



Then $\left(\theta_{1}, \theta_{2}\right)=(X Y, Y)$ is a transformation with the Jacobian $|J|=\theta_{2}^{-1}$ and the density written as follows,

$$
\begin{aligned}
f_{\theta_{1}, \theta_{2}}\left(\theta_{1}, \theta_{2}\right)= & f_{X, Y}\left(\frac{\theta_{1}}{\theta_{2}}, \theta_{2}\right) \theta_{2}^{-1} \\
& =\frac{S / 2}{\Gamma(1 / 2)}\left(\frac{\theta_{1}}{\theta_{2}}\right)^{1 / 2-1} e^{-\frac{\theta_{1}}{\theta_{2}} S / 2} \\
& \frac{S / 2}{\Gamma((J+2) / 2)} \theta_{2}^{-(J+2) / 2-1} e^{-\frac{S / 2}{\theta_{2}}} \theta_{2}^{-1}
\end{aligned}
$$

Next $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ is proportional to $f_{\theta_{1}, \theta_{2}}\left(\theta_{1}, \theta_{2}\right)$ and so we can derive the conditional density of $\theta_{2}$ as follows,

$$
\begin{align*}
f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right) & \propto f_{\theta_{1}, \theta_{2}\left(\theta_{1}, \theta_{2}\right)}^{(\mathrm{C} 10}  \tag{C10}\\
& \propto \theta_{2}^{1-1 / 2} e^{-\frac{1}{\theta_{2}} \theta_{1} S / 2} \theta_{2}^{-(J+2) / 2-1} e^{-\frac{1}{\theta_{2}} S / 2} \theta_{2}^{-1}  \tag{C11}\\
& =\theta_{2}^{-(1 / 2+(J+2) / 2)-1} e^{-\frac{1}{\theta_{2}}\left(\theta_{1}+1\right) S / 2}  \tag{C12}\\
& =\theta_{2}^{-(J-1) / 2-1} e^{-\frac{1}{\theta_{2}}\left(\theta_{1}+1\right) S / 2} \tag{C13}
\end{align*}
$$

This is proportional to an inverse gamma distribution and so, $\theta_{2} \left\lvert\, \theta_{1} \sim \Gamma^{-1}\left(\frac{J-1}{2},\left(\theta_{1}+1\right) S / 2\right)\right.$. We know that the inverse gamma distribution is continuous and unimodal, so $M=1$.
(b) Calculate the maximum value of $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ : Similar to Fig. 5 of Example 4.7, $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ reaches its maximum height when $\theta_{1}=0$. It can also be shown from equation C 10 that the density function of $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ is maximized when $\theta_{1}=0$ since the normalizing constant will be the largest. So the largest value of $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}\right)$ will occur when $\theta_{1}=0$. To find the maximum conditional distribution, we find the value of $f_{\theta_{2} \mid \theta_{1}}\left(\theta_{2} \mid \theta_{1}=0\right)$ evaluated
at $\theta_{2}=\frac{S}{J+1}$, the mode (see Section 5.3 of Hoff (2009)).

$$
\begin{aligned}
K & =f_{\theta_{2} \mid \theta_{1}}\left(\left.\frac{S}{J+1} \right\rvert\, \theta_{1}=0\right) \\
& =\left.\frac{(S / 2)^{\frac{J-1}{2}}}{\Gamma\left(\frac{J-1}{2}\right)} y^{-\frac{J-1}{2}-1} e^{-\frac{S / 2}{y}}\right|_{y=\frac{S}{J+1}} \\
& =\frac{(S / 2)^{\frac{J-1}{2}}}{\Gamma\left(\frac{J-1}{2}\right)}\left(\frac{S}{J+1}\right)^{-\frac{J-3}{2}} e^{-\frac{J+1}{2}}
\end{aligned}
$$

And so,
$K=\frac{(S / 2)^{\frac{J-1}{2}}}{\Gamma\left(\frac{J-1}{2}\right)}\left(\frac{S}{J+1}\right)^{-\frac{J-3}{2}} e^{-\frac{J+1}{2}}$

Proof of Lemma 3.5 By the property of stationary distribution, if $\sigma_{n-1}^{2} \sim \pi$ then $\sigma_{n}^{2} \sim \pi$ and so the lemma follows from the following.

$$
\begin{aligned}
E_{\sigma_{n}^{2} \sim \pi}\left[V\left(\sigma_{n}^{2}\right)\right] & =E_{\sigma_{n-1}^{2} \sim \pi}\left[E\left[V\left(\sigma_{n}^{2}\right) \mid \sigma_{n-1}^{2}\right]\right] \\
& \leq E_{\sigma_{n-1}^{2} \sim \pi}\left[\lambda V\left(\sigma_{n-1}^{2}\right)+b\right] \\
& =\lambda E_{\sigma_{n}^{2} \sim \pi}\left[V\left(\sigma_{n}^{2}\right)\right]+b
\end{aligned}
$$

Proof of 4.19 Let $\lambda=0.6583702, h=-0.5248723$ and $b=106.3874$, then

$$
\begin{aligned}
& E\left[V\left(\sigma_{n}^{2}\right) \mid \sigma_{n-1}^{2}\right] \\
& =E\left[\left(\sigma_{n}^{2}-h\right)^{2} \mid \sigma_{n-1}^{2}\right] \\
& =E\left[\left(\sigma_{n}^{2}\right)^{2}-2 h \sigma_{n}^{2}+h^{2} \mid \sigma_{n-1}^{2}\right] \\
& =E\left[\left(X_{n} Y_{n} \sigma_{n-1}^{2}+Y_{n}\right)^{2}-2 h\left(X_{n} Y_{n} \sigma_{n-1}^{2}+Y_{n}\right)+h^{2} \mid \sigma_{n-1}^{2}\right]
\end{aligned}
$$

$$
\begin{aligned}
= & E\left[Y_{n}^{2}\right]\left(E\left[X_{n}^{2}\right]\left(\sigma_{n-1}^{2}\right)^{2}+2 E\left[X_{n}\right] \sigma_{n-1}^{2}+1\right) \\
& -2 h\left(E\left[X_{n}\right] E\left[Y_{n}\right] \sigma_{n-1}^{2}+E\left[Y_{n}\right]\right)+h^{2} \\
= & E\left[Y_{n}^{2}\right] E\left[X_{n}^{2}\right]\left(\sigma_{n-1}^{2}\right)^{2}+2 E\left[X_{n}\right] E\left[Y_{n}^{2}\right] \sigma_{n-1}^{2} \\
& +E\left[Y_{n}^{2}\right]-2 h E\left[X_{n}\right] E\left[Y_{n}\right] \sigma_{n-1}^{2}-2 h E\left[Y_{n}\right]+h^{2} \\
= & E\left[Y_{n}^{2}\right] E\left[X_{n}^{2}\right]\left(\sigma_{n-1}^{2}\right)^{2}+2 E\left[X_{n}\right]\left(E\left[Y_{n}^{2}\right]\right. \\
& \left.-h E\left[Y_{n}\right]\right) \sigma_{n-1}^{2}+E\left[Y_{n}^{2}\right]-2 h E\left[Y_{n}\right]+h^{2} \\
= & 0.6583702\left(\sigma_{n-1}^{2}\right)^{2}+0.6911206 \sigma_{n-1}^{2}+107.3691 \\
= & \lambda\left(\sigma_{n-1}^{2}\right)^{2}+2 \lambda h \sigma_{n-1}^{2}+\lambda h^{2}+b \\
= & \lambda\left(\sigma_{n-1}^{2}+h\right)^{2}+b
\end{aligned}
$$

## C. 4 Proof of Theorem 4.23

Proof of Theorem 4.23 This example uses a modified version of the Sideways Theorem 4.2 to find an upper bound on the convergence rate. We will also use Proposition 2.2, which states that the total variation between two random variables is equal to the total variation of any invertible transformation of the same two random variables.

Let $\vec{X}_{n}, \vec{X}_{n}^{\prime} \in \mathbb{R}^{2}$ be two copies of the autoregressive normal process as defined in Example 4.22. Then for $\vec{Z}_{n} \sim$ $N\left(\overrightarrow{0}, I_{d}\right)$,

$$
\vec{X}_{n}=A \vec{X}_{n-1}+\Sigma_{d} \vec{Z}_{n} \quad \vec{X}_{n}^{\prime}=A \vec{X}_{n-1}^{\prime}+\Sigma_{d} \vec{Z}_{n}^{\prime}
$$

We apply the one-shot coupling method to bound the total variation distance. For $n<N$ set $\vec{Z}_{n}=\vec{Z}_{n}^{\prime}$.

Suppose $X_{0}, X_{0}^{\prime}$ are known and define

$$
\Delta=\left\|\Sigma_{d}^{-1} A^{n}\left(\vec{X}_{0}-\vec{X}_{0}^{\prime}\right)\right\|_{2}
$$

Decompose $A=P D P^{-1}$ with $D$ as the corresponding diagonal matrix, $\lambda_{i}$ is the $i$ th eigenvalue of $A$ and $\|\cdot\|_{2}$ denotes the Frobenius norm. Then $\Delta$ is bounded above as follows,

$$
\begin{aligned}
\Delta & =\left\|\Sigma_{d}^{-1} A^{n}\left(\vec{X}_{0}-\vec{X}_{0}^{\prime}\right)\right\|_{2} \\
& =\left\|\Sigma_{d}^{-1} P D^{n} P^{-1}\left(\vec{X}_{0}-\vec{X}_{0}^{\prime}\right)\right\|_{2} \\
& \leq\left\|\Sigma_{d}^{-1}\right\|_{2} \cdot\left\|\left.P\right|_{2}\right\| D^{n}\left\|_{2}\right\| P^{-1}\left\|_{2}\right\| \vec{X}_{0}-\vec{X}_{0}^{\prime} \|_{2}
\end{aligned}
$$

by Lemma 1.2.7 of Aggarwal (2020)

$$
\begin{aligned}
& \leq\left\|\Sigma_{d}^{-1}\right\|_{2} \cdot\|P\|_{2}\left\|P^{-1}\right\|_{2}\left\|\vec{X}_{0}-\vec{X}_{0}^{\prime}\right\|_{2} \sqrt{\sum_{i=1}^{d}\left|\lambda_{i}\right|^{2 n}} \\
& \leq\left\|\Sigma_{d}^{-1}\right\|_{2} \cdot\|P\|_{2}\left\|P^{-1}\right\|_{2}\left\|\vec{X}_{0}-\vec{X}_{0}^{\prime}\right\|_{2} \sqrt{d} \max _{1 \leq i \leq d}\left|\lambda_{i}\right|^{n}
\end{aligned}
$$

For now assume that $X_{0}, X_{0}^{\prime}$ are known and note that $\Sigma_{d}^{-1}$ is an invertible transform. We bound the total variation distance as follows by applying two invertible transforms on the Markov chain and using the fact that $\vec{Z}_{m}=\vec{Z}_{m}^{\prime}, m<N$.

$$
\begin{aligned}
& \left\|\mathcal{L}\left(\vec{X}_{N}\right)-\mathcal{L}\left(\vec{X}_{N}^{\prime}\right)\right\| \\
& \leq E_{\left\{\vec{Z}_{m}\right\}_{m<N}}\left[\left\|\mathcal{L}\left(\vec{X}_{N}\right)-\mathcal{L}\left(\vec{X}_{N}^{\prime}\right)\right\|\right]
\end{aligned}
$$

by Proposition 2.3
$=E_{\left\{\vec{Z}_{m}\right\}_{m<N}}\left[\left\|\mathcal{L}\left(\Sigma_{d}^{-1} \vec{X}_{N}\right)-\mathcal{L}\left(\Sigma_{d}^{-1} \vec{X}_{N}^{\prime}\right)\right\|\right]$
by Proposition 2.2
$=E_{\left\{\vec{Z}_{m}\right\}_{m<N}}\left[\left\|\mathcal{L}\left(\Sigma_{d}^{-1} A \vec{X}_{N-1}+\vec{Z}_{N}\right)-\mathcal{L}\left(\Sigma_{d}^{-1} A \vec{X}_{N-1}^{\prime}+\vec{Z}_{N}^{\prime}\right)\right\|\right]$
$=E_{\left\{\vec{Z}_{m}\right\}_{m<N}}\left[\left\|\mathcal{L}\left(\Sigma_{d}^{-1} A^{N} \vec{X}_{0}+\vec{Z}_{N}\right)-\mathcal{L}\left(\Sigma_{d}^{-1} A^{N} \vec{X}_{0}^{\prime}+\vec{Z}_{N}^{\prime}\right)\right\|\right]$
by Proposition 2.2

$$
\begin{aligned}
& =E_{\left\{\vec{Z}_{m}\right\}_{m<N}}\left[\left\|\mathcal{L}\left(\vec{Z}_{N}+\Sigma_{d}^{-1} A^{N}\left(\vec{X}_{0}-\vec{X}_{0}^{\prime}\right)\right)-\mathcal{L}\left(\vec{Z}_{N}^{\prime}\right)\right\|\right] \\
& =\left\|\mathcal{L}\left(\vec{Z}_{N}+\Sigma_{d}^{-1} A^{N}\left(\vec{X}_{0}-\vec{X}_{0}^{\prime}\right)\right)-\mathcal{L}\left(\vec{Z}_{N}^{\prime}\right)\right\|
\end{aligned}
$$

There exists a rotation matrix $R \in \mathbb{R}^{d \times d}$ such that

$$
\begin{aligned}
R\left[\Sigma_{d}^{-1} A\left(\vec{X}_{n}-\vec{X}_{n}^{\prime}\right)\right] & =\left(\left\|\Sigma_{d}^{-1} A\left(\vec{X}_{n}-\vec{X}_{n}^{\prime}\right)\right\|_{2}, 0, \ldots 0\right) \\
& =(\Delta, 0, \ldots 0)
\end{aligned}
$$

Aggarwal (2020). By properties of rotation, $R$ is orthogonal, so $R^{T}=R^{-1}$ and $R Z_{n} \sim N\left(0, R I_{d} R^{T}\right)=N\left(0, I_{d}\right) \sim Z_{n}$. In other words, $R Z_{n} \stackrel{d}{=} Z_{n} \stackrel{d}{=} Z_{n}^{\prime}$. Thus, continuing the above equality,

$$
\begin{aligned}
& \left\|\mathcal{L}\left(\vec{X}_{n}\right)-\mathcal{L}\left(\vec{X}_{n}^{\prime}\right)\right\| \\
& \quad \leq\left\|\mathcal{L}\left(\vec{Z}_{n}+\Sigma_{d}^{-1} A^{n}\left(\vec{X}_{0}-\vec{X}_{0}^{\prime}\right)\right)-\mathcal{L}\left(\vec{Z}_{n}^{\prime}\right)\right\| \\
& \quad=\left\|\mathcal{L}\left(R\left[\vec{Z}_{n}+\Sigma_{d}^{-1} A\left(\vec{X}_{n}-\vec{X}_{n}^{\prime}\right)\right]\right)-\mathcal{L}\left(R \vec{Z}_{n}^{\prime}\right)\right\| \quad \text { by Proposition } 2.2 \\
& \quad=\left\|\mathcal{L}\left(\vec{Z}_{n}+(\Delta, 0, \ldots 0)\right)-\mathcal{L}\left(\vec{Z}_{n}\right)\right\|
\end{aligned}
$$

Next, suppose that $X_{0}, X_{0}^{\prime}$ are unknown. Then, the inequality stated in Eq. 12 is shown as follows,

$$
\begin{aligned}
& \left\|\mathcal{L}\left(\vec{X}_{n}\right)-\mathcal{L}\left(\vec{X}_{n}^{\prime}\right)\right\| \\
& \leq E_{\Delta}\left[\left\|\mathcal{L}\left(\vec{Z}_{n}+(\Delta, 0, \ldots 0)\right)-\mathcal{L}\left(\vec{Z}_{n}\right)\right\|\right] \quad \text { by Proposition } 2.3 \\
& =E_{\Delta}\left[\frac{1}{2} \int_{\mathbb{R}^{d}} \frac{e^{-\sum_{i=2}^{d} y_{i}^{2} / 2}}{(2 \pi)^{d / 2}}\left|e^{-y_{1}^{2} / 2}-e^{-\left(y_{1}-\Delta\right)^{2} / 2}\right| d \vec{y}\right] \\
& =E_{\Delta}\left[\frac{1}{2} \int_{\mathbb{R}}\left|\frac{1}{\sqrt{2 \pi}} e^{-y_{1}^{2} / 2}-\frac{1}{\sqrt{2 \pi}} e^{-\left(y_{1}-\Delta\right)^{2} / 2} d\right| \vec{y}\right] \\
& =E_{\Delta}\left[\left\|\mathcal{L}\left(Z_{1, n}+\Delta\right)-\mathcal{L}\left(Z_{1, n}\right)\right\|\right] \\
& \leq \frac{1}{\sqrt{2 \pi}} E[\Delta] \text { by Lemma } B .3 \\
& \leq \sqrt{\frac{d}{2 \pi}}\left\|\Sigma_{d}^{-1}\right\|_{2} \cdot\|P\|_{2}\left\|P^{-1}\right\|_{2} E\left[\left\|\vec{X}_{0}-\vec{X}_{0}^{\prime}\right\|_{2}\right] \max _{1 \leq i \leq d}\left|\lambda_{i}\right|^{n}
\end{aligned}
$$

## Appendix D: Lemmas for ARCH process examples

## D. 1 Proof of lemmas used in Theorem 5.3

Proof of Lemma 5.4 Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1} \in \mathbb{R}$ be two copies of the LARCH process. For fixed $n \geq 1$, let $Z_{n}=Z_{n}^{\prime}$ and so,

$$
\begin{aligned}
E\left[\left|X_{n}-X_{n}^{\prime}\right|\right] & =E\left[\left|\left(\beta_{0}+\beta_{1} X_{n-1}\right) Z_{n}-\left(\beta_{0}+\beta_{1} X_{n-1}^{\prime}\right) Z_{n}\right|\right] \\
& \leq \beta_{1} E\left[\left|Z_{n}\right|\right] E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]
\end{aligned}
$$

Since $Z_{n} \stackrel{d}{=} Z_{0}>0$ a.s., the geometric convergence rate is $D=\beta_{1} E\left[Z_{0}\right]$.

Proof of Lemma 5.5 For a fixed $n \geq 0$, suppose that $Z_{n+1}$, $Z_{n+1}^{\prime}$ are independent. By Proposition 2.3, the total variation distance between the two processes is bounded above by the expectation of the total variation.

$$
\begin{aligned}
\left\|\mathcal{L}\left(X_{n+1}\right)-\mathcal{L}\left(X_{n+1}^{\prime}\right)\right\| \leq & E\left[\| \mathcal{L}\left(\left(\beta_{0}+\beta_{1} X_{n}\right) Z_{n+1}\right)\right. \\
& \left.-\mathcal{L}\left(\left(\beta_{0}+\beta_{1} X_{n}^{\prime}\right) Z_{n+1}\right) \|\right]
\end{aligned}
$$

Note that $Z_{n+1}$ and $Z_{n+1}^{\prime}$ are used interchangeably in the total variation distance since $Z_{n+1} \stackrel{d}{=} Z_{n+1}^{\prime}$. Let $Y_{n}=\beta_{0}+\beta_{1} X_{n}$, $Y_{n}^{\prime}=\beta_{0}+\beta_{1} X_{n}^{\prime}, \Delta=Y_{n}^{\prime}-Y_{n}$, and $\Delta^{\prime}=\frac{\Delta}{Y_{n}}$. WLOG $Y_{n}^{\prime}>Y_{n}$ so that $\Delta, \Delta^{\prime}>0$. Then,

```
\(\left\|\mathcal{L}\left(X_{n+1}\right)-\mathcal{L}\left(X_{n+1}^{\prime}\right)\right\|\)
    \(\leq E\left[\left\|\mathcal{L}\left(Y_{n} Z_{n+1}\right)-\mathcal{L}\left(Y_{n}^{\prime} Z_{n+1}\right)\right\|\right]\) by Proposition 2.3
    \(=E\left[\left\|\mathcal{L}\left(Y_{n} Z_{n+1}\right)-\mathcal{L}\left(\left(Y_{n}+\Delta\right) Z_{n+1}\right)\right\|\right]\)
    \(=E\left[\left\|\mathcal{L}\left(Z_{n+1}\right)-\mathcal{L}\left(\left(1+\Delta^{\prime}\right) Z_{n+1}\right)\right\|\right]\) by Proposition 2.2
    \(=E\left[\left\|\mathcal{L}\left(\log \left(Z_{n+1}\right)\right)-\mathcal{L}\left(\log \left(1+\Delta^{\prime}\right)+\log \left(Z_{n+1}\right)\right)\right\|\right]\)
```

        by Proposition 2.2
    $$
\begin{aligned}
& \leq \frac{M+1}{2} \sup _{x} e^{x} f_{Z_{n}}\left(e^{x}\right) E\left[\log \left(1+\Delta^{\prime}\right)\right] \\
& \leq \frac{M+1}{2} \sup _{x} e^{x} f_{Z_{n}}\left(e^{x}\right) \frac{E[|\Delta|]}{\beta_{0}} \\
& =\frac{M+1}{2} \sup _{x} e^{x} f_{Z_{n}}\left(e^{x}\right) \frac{\beta_{1} E\left[\left|X_{n}-X_{n}^{\prime}\right|\right]}{\beta_{0}}
\end{aligned}
$$

The second last inequality is by Lemma B.3. See the proof of Lemma 4.3 for more details. The last inequality is by the Mean Value Theorem.

## D. 2 Proof of lemmas used in Theorem 5.8

Proof of Lemma 5.9 Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1} \in \mathbb{R}$ be two copies of the asymmetric ARCH process.

For a fixed $n \geq 1$, let $Z_{n}=Z_{n}^{\prime}$ and so,
$E\left[\left|X_{n}-X_{n}^{\prime}\right|\right]$

$$
\begin{aligned}
&= E\left[\left|\sqrt{\left(a X_{n-1}+b\right)^{2}+c^{2}} Z_{n}-\sqrt{\left(a X_{n-1}^{\prime}+b\right)^{2}+c^{2}} Z_{n}\right|\right. \\
&=\left|\sqrt{\left(a X_{n-1}+b\right)^{2}+c^{2}}-\sqrt{\left(a X_{n-1}^{\prime}+b\right)^{2}+c^{2}}\right| \\
& E\left[\left|Z_{n}\right|\right]
\end{aligned}
$$

Note that the derivative of $f(x)=\sqrt{(a x+b)^{2}+c^{2}}$ is

$$
\begin{equation*}
\left|f^{\prime}(x)\right|=\left|\frac{a(a x+b)}{\sqrt{(a x+b)^{2}+c^{2}}}\right| \leq \frac{|a(a x+b)|}{\sqrt{(a x+b)^{2}}}=|a| \tag{D15}
\end{equation*}
$$

and so,
$E\left[\left|X_{n}-X_{n}^{\prime}\right|\right] \leq|a| E\left[\left|Z_{n}\right|\right] E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]$
Thus, the geometric convergence rate is $D=|a| E\left[\left|Z_{0}\right|\right]$.
Proof of Lemma 5.10 Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1} \in \mathbb{R}$ be two copies of the asymmetric ARCH process.

For $n \geq 1, Z_{n}, Z_{n}^{\prime}$ are independent. By Proposition 2.3, the total variation distance between the two processes is bounded above by the expectation of the total variation with respect to $X_{n-1}, X_{n-1}^{\prime}, Z_{n}, Z_{n}^{\prime}$.

$$
\begin{aligned}
& \left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \\
& \quad \leq E\left[\| \mathcal{L}\left(\sqrt{\left(a X_{n-1}+b\right)^{2}+c^{2}} Z_{n}\right)\right. \\
& \left.\quad-\mathcal{L}\left(\sqrt{\left(a X_{n-1}^{\prime}+b\right)^{2}+c^{2}} Z_{n}^{\prime}\right) \|\right]
\end{aligned}
$$

Let $Y_{n-1}=\sqrt{\left(a X_{n-1}+b\right)^{2}+c^{2}}$ and $Y_{n-1}^{\prime}$ $=\sqrt{\left(a X_{n-1}+b\right)^{2}+c^{2}}, \Delta=Y_{n-1}^{\prime}-Y_{n-1}$ and $\Delta^{\prime}=$ $\frac{\Delta}{Y_{n-1}}$. WLOG, $Y_{n-1}^{\prime}<Y_{n-1}$, so $-1<\Delta^{\prime}<0$, because $Y_{n-1}, Y_{n-1}^{\prime}>0$ and

$$
\begin{aligned}
& \left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq E\left[\left\|\mathcal{L}\left(Y_{n-1} Z_{n}\right)-\mathcal{L}\left(Y_{n-1}^{\prime} Z_{n}\right)\right\|\right] \\
& \quad=E\left[\left\|\mathcal{L}\left(Y_{n-1} Z_{n}\right)-\mathcal{L}\left(\left(Y_{n-1}+\Delta\right) Z_{n}\right)\right\|\right]
\end{aligned}
$$

by Proposition 2.2

$$
\begin{aligned}
& =E\left[\left\|\mathcal{L}\left(Z_{n}\right)-\mathcal{L}\left(\left(1+\Delta^{\prime}\right) Z_{n}\right)\right\|\right] \quad \text { by Proposition } 2.2 \\
& \leq E\left[\sup _{x} 1-\frac{\pi_{Z_{n}}(x)}{\pi_{\left(1+\Delta^{\prime}\right) Z_{n}}(x)}\right]
\end{aligned}
$$

by Lemma 6.16 of Levinet al. (2017)
Let the density of $Z_{n}$ be $\pi_{Z_{n}}(x)$, then $\pi_{\left(1+\Delta^{\prime}\right) Z_{n}}(x)=$ $\frac{1}{1+\Delta^{\prime}} \pi_{Z_{n}}\left(\frac{x}{1+\Delta^{\prime}}\right)$.

$$
\begin{aligned}
& \left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq E\left[\sup _{x} 1-\left(1+\Delta^{\prime}\right) \frac{\pi_{Z_{n}}(x)}{\pi_{Z_{n}}\left(\frac{x}{1+\Delta^{\prime}}\right)}\right] \\
& \quad \leq E\left[\sup _{x} 1-\left(1+\Delta^{\prime}\right)\right] \\
& \quad=E\left[\Delta^{\prime}\right]
\end{aligned}
$$

$$
\begin{aligned}
& \leq \frac{E\left[\left|Y_{n-1}-Y_{n-1}^{\prime}\right|\right]}{c} \text { since } Y_{n-1} \geq c \\
& \leq \frac{|a|}{c} E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]
\end{aligned}
$$

by equation $D 15$
The second inequality is by assumption $\pi_{Z_{n}}(x) \geq \pi_{Z_{n}}\left(\frac{x}{1+\Delta^{\prime}}\right)$.

## D. 3 Proof of lemmas used in Theorem 5.13

Proof of Lemma 5.14 Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1} \in \mathbb{R}$ be two copies of the GARCH process. For $n \geq 2$, let $Z_{n}=Z_{n}^{\prime}$. First note that,

$$
\begin{align*}
E\left[\left|X_{n}-X_{n}^{\prime}\right|\right] & =E\left[\left|\sigma_{n} Z_{n}-\sigma_{n}^{\prime} Z_{n}\right|\right]=E\left[\left|\sigma_{n}-\sigma_{n}^{\prime}\right|\left|Z_{n}\right|\right] \\
& =E\left[\left|\sigma_{n}-\sigma_{n}^{\prime}\right|\right] E\left[\left|Z_{n}\right|\right] \tag{D16}
\end{align*}
$$

Next, we find an upper bound on $E\left[\left|\sigma_{n}-\sigma_{n}^{\prime}\right|\right]$ by first noting that $\sigma_{n}^{2}=\alpha^{2}+\left(\beta^{2} Z_{n-1}^{2}+\gamma^{2}\right) \sigma_{n-1}^{2}$ by substitution.

$$
\begin{aligned}
& E\left[\left|\sigma_{n}-\sigma_{n}^{\prime}\right|\right] \\
& \quad=E\left[\mid \sqrt{\alpha^{2}+\left(\beta^{2} Z_{n-1}^{2}+\gamma^{2}\right) \sigma_{n-1}^{2}}\right. \\
& \left.\quad-\sqrt{\alpha^{2}+\left(\beta^{2} Z_{n-1}^{2}+\gamma^{2}\right) \sigma_{n-1}^{\prime 2} \mid}\right] \\
& \quad \leq E\left[\sqrt{\beta^{2} Z_{n-1}^{2}+\gamma^{2}}\right] E\left[\left|\sigma_{n-1}-\sigma_{n-1}^{\prime}\right|\right] \\
& \quad=E\left[\sqrt{\beta^{2} Z_{n-1}^{2}+\gamma^{2}}\right] \frac{E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]}{E\left[\left|Z_{n-1}\right|\right]}
\end{aligned}
$$

The above inequality is by taking the maximum of the derivative and the last equality is a result of Eq. D16. Finally, substituting $E\left[\left|\sigma_{n}-\sigma_{n}^{\prime}\right|\right]$ into Eq. D16,

$$
\begin{aligned}
& E\left[\left|X_{n}-X_{n}^{\prime}\right|\right] \\
& \quad \leq E\left[\sqrt{\beta^{2} Z_{n-1}^{2}+\gamma^{2}}\right] \frac{E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]}{E\left[\left|Z_{n-1}\right|\right]} E\left[\left|Z_{n}\right|\right] \\
& \quad=E\left[\sqrt{\beta^{2} Z_{n-1}^{2}+\gamma^{2}}\right] E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right] \\
& \\
& \quad \leq \sqrt{\beta^{2} E\left[Z_{0}^{2}\right]+\gamma^{2}} E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]
\end{aligned}
$$

by Jensen's inequality

Thus, the geometric convergence rate is
$D=\sqrt{\beta^{2} E\left[Z_{0}^{2}\right]+\gamma^{2}}$.
Proof of Lemma 5.15 Let $\left\{X_{n}\right\}_{n \geq 1} \in \mathbb{R}$ and $\left\{X_{n}^{\prime}\right\}_{n \geq 1} \in \mathbb{R}$ be two copies of the GARCH process.

For $n \geq 2$, suppose that $Z_{n}, Z_{n}^{\prime}$ are independent. By Proposition 2.3, the total variation distance between the two
processes is bounded above by the expectation of the total variation.

$$
\left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \leq E\left[\left\|\mathcal{L}\left(\sigma_{n} Z_{n}\right)-\mathcal{L}\left(\sigma_{n}^{\prime} Z_{n}\right)\right\|\right]
$$

Let $\Delta=\sigma_{n}^{\prime}-\sigma_{n}$ and $\Delta^{\prime}=\frac{\Delta}{\sigma_{n}}$. WLOG, $\sigma_{n}^{\prime}<\sigma_{n}$, so $\Delta, \Delta^{\prime}<$ 0 because $\sigma_{n}, \sigma_{n}^{\prime}>0$ and

$$
\begin{aligned}
& \left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \\
& \quad=E\left[\left\|\mathcal{L}\left(\sigma_{n} Z_{n}\right)-\mathcal{L}\left(\left(\sigma_{n}+\Delta\right) Z_{n}\right)\right\|\right]
\end{aligned}
$$

by Proposition 2.2

$$
=E\left[\left\|\mathcal{L}\left(Z_{n}\right)-\mathcal{L}\left(\left(1+\Delta^{\prime}\right) Z_{n}\right)\right\|\right]
$$

by Proposition 2.2

$$
\leq E\left[\sup _{x} 1-\frac{\pi_{Z_{n}}(x)}{\pi_{\left(1+\Delta^{\prime}\right) Z_{n}}(x)}\right]
$$

by Lemma 6.16 of Levinet al. (2017)
Let the density of $Z_{n}$ be $\pi_{Z_{n}}(x)$, then $\pi_{\left(1+\Delta^{\prime}\right) Z_{n}}(x)=$ $\frac{1}{1+\Delta^{\prime}} \pi_{Z_{n}}\left(\frac{x}{1+\Delta^{\prime}}\right)$.

$$
\begin{aligned}
& \left\|\mathcal{L}\left(X_{n}\right)-\mathcal{L}\left(X_{n}^{\prime}\right)\right\| \\
& \quad \leq E\left[\sup _{x} 1-\left(1+\Delta^{\prime}\right) \frac{\pi_{Z_{n}}(x)}{\pi_{Z_{n}}\left(\frac{x}{1+\Delta^{\prime}}\right)}\right] \\
& \quad \leq E\left[\sup _{x} 1-\left(1+\Delta^{\prime}\right)\right]
\end{aligned}
$$

$$
\text { by assumption } \pi_{Z_{n}}(x) \geq \pi_{Z_{n}}\left(\frac{x}{1+\Delta^{\prime}}\right)
$$

$$
=E\left[\Delta^{\prime}\right]
$$

$$
\leq \frac{E\left[\left|\sigma_{n}^{\prime}-\sigma_{n}\right|\right]}{\alpha} \text { since } \sigma_{n} \geq \alpha
$$

$$
\leq \frac{D}{\alpha E\left[\left|Z_{n-1}\right|\right]} E\left[\left|X_{n-1}-X_{n-1}^{\prime}\right|\right]
$$

by equation in proof $D .3$

Proof of Lemma 5.16

$$
\begin{aligned}
& E\left[\left|X_{1}-X_{1}^{\prime}\right|\right] \\
& =\left|\sigma_{1}^{2}-\sigma_{1}^{\prime 2}\right| E\left[\left|Z_{1}\right|\right] \quad \text { by equation in proof } D .3 \\
& =\left|\sqrt{\alpha^{2}+\beta^{2} X_{0}^{2}+\gamma^{2} \sigma_{0}^{2}}-\sqrt{\alpha^{2}+\beta^{2} X_{0}^{\prime 2}+\gamma^{2} \sigma_{0}^{\prime 2}}\right| E\left[\left|Z_{1}\right|\right] \\
& \leq \sqrt{\left|\left(\alpha^{2}+\beta^{2} X_{0}^{2}+\gamma^{2} \sigma_{0}^{2}\right)-\left(\alpha^{2}+\beta^{2} X_{0}^{\prime 2}+\gamma^{2} \sigma_{0}^{\prime 2}\right)\right| E\left[\left|Z_{1}\right|\right]} \\
& \text { since }|\sqrt{x}-\sqrt{y}|=\sqrt{(\sqrt{x}-\sqrt{y})^{2}}=\sqrt{x+y-2 \sqrt{x} \sqrt{y}} \\
& \leq \sqrt{|x-y|} \\
& \leq \sqrt{\beta^{2}\left|X_{0}^{2}-X_{0}^{\prime 2}\right|+\gamma^{2}\left|\sigma_{0}^{2}-\sigma_{0}^{\prime 2}\right| E\left[\left|Z_{0}\right|\right]}
\end{aligned}
$$

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[^0]:    Sabrina Sixta
    sabrina.sixta@mail.utoronto.ca
    Jeffrey S. Rosenthal
    jeff@math.toronto.edu
    1 Department of Statistical Sciences, University of Toronto, 700 University Avenue, 9th Floor, Toronto, ON M5G 1Z5, Canada

[^1]:    ${ }^{1}$ The trigonometric identities used are $2 \cos \mu \sin v=\sin (\mu+v)-$ $\sin (\mu-v)$ and $\cos (\mu+v)=\cos \mu \cos v+\sin \mu \sin v$ where $\mu, v \in \mathbb{R}$

