

# Bootstrapping non-stationary stochastic volatility

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## A Online Appendix

### A.1 Auxiliary results

Throughout, we make use of the following version of Skorokhod's representation theorem.

**Theorem A.1.** [Kallenberg, 1997, Corollary 5.12] *Let  $f$  and  $\{f_n\}_{n \geq 1}$  be measurable functions from a Borel space  $\mathcal{S}$  to a Polish space  $\mathcal{T}$ , and let  $\xi$  and  $\{\xi_n\}_{n \geq 1}$  be random elements in  $\mathcal{S}$  with  $f_n(\xi_n) \xrightarrow{w} f(\xi)$ . Then there exist some random elements  $\tilde{\xi} \stackrel{d}{=} \xi$  and  $\tilde{\xi}_n \stackrel{d}{=} \xi_n$  defined on a common probability space with  $f_n(\tilde{\xi}_n) \xrightarrow{a.s.} f(\tilde{\xi})$ .*

The next lemma contains a result about the asymptotic continuity of the distribution function of Dickey-Fuller type-statistics under non-stationary stochastic volatility.

**Lemma A.1.** *With  $M$  and  $V$  defined in Lemma 1, under Assumptions 1 and 2, let*

$$\tau_1 := \frac{\int_0^1 M(u) dM(u)}{\int_0^1 M^2(u) du} \quad \text{and} \quad \tau_2 := \frac{\int_0^1 M(u) dM(u)}{\sqrt{V(1) \int_0^1 M^2(u) du}}.$$

*Then the random cdfs  $F_1(\cdot) := P(\tau_1 \leq \cdot | \sigma)$  and  $F_2(\cdot) := P(\tau_2 \leq \cdot | \sigma)$  are sample-path continuous a.s.*

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PROOF OF LEMMA A.1. We reduce the proof to the following well-known result (a, a, pp. 472–473). Let  $\{X(u)\}_{u \in [0,1]}$  be a Gaussian process with mean zero and a continuous covariance kernel, let  $q : [0, 1] \rightarrow \mathbb{R}$  be a square-integrable function and let  $\alpha \in \mathbb{R}$  be arbitrary. Then the distribution of  $\int_0^1 (X(u) + \alpha q(u))^2 du$  is that of an infinite series of independent non-central  $\chi^2$  random variables and, as a result, it has a continuous cdf.

The random cdfs  $F_1$  and  $F_2$  are determined, up to a modification, by the *distribution* of  $(B_z, \sigma)$ , such that the structure of the probability space on which  $(B_z, \sigma)$  is defined is irrelevant for the claim of interest. We therefore assume, without loss of generality, that the independent processes  $B_z$  and  $\sigma$  are defined on a product probability space. Let  $(\Omega_\sigma, \mathcal{F}_\sigma, P_\sigma)$  be the factor-space on which  $\sigma$  is defined. Fix  $A \in \mathcal{F}_\sigma$  with  $P_\sigma(A) = 1$  such that  $V(\omega, \cdot) := \int_0^1 \sigma^2(\omega, u) du$  is well-defined, continuous and  $0 < V(\omega, 1) < \infty$ . Let  $\Gamma := \{\sigma(\omega, \cdot) : \omega \in A\}$  be the set of trajectories for  $\sigma$  when  $\omega \in A$ . For every  $\gamma \in \Gamma$ , the process  $M_\gamma(\cdot) := \int_0^\cdot \gamma(u) dB_z(u)$  is a.s. well-defined and  $\int_0^1 M_\gamma^2(u) du > 0$  a.s. The result in the lemma will follow if the deterministic cdfs  $P(\tau_{\gamma 1} \leq \cdot)$  and  $P(\tau_{\gamma 2} \leq \cdot)$  are continuous for every  $\gamma \in \Gamma$ :

$$P(\tau_{\gamma 1} = x) = 0, \quad P(\tau_{\gamma 2} = x) = 0, \quad \forall (x, \gamma) \in \mathbb{R} \times \Gamma, \quad (\text{A.1})$$

where

$$\tau_{\gamma 1} := \frac{\int_0^1 M_\gamma(u) dM_\gamma(u)}{\int_0^1 M_\gamma^2(u) du}, \quad \tau_{\gamma 2} := \frac{\int_0^1 M_\gamma(u) dM_\gamma(u)}{\sqrt{V(1) \int_0^1 M_\gamma^2(u) du}}.$$

In fact, (A.1) implies that  $F_1$  and  $F_2$  have sample-path continuous modifications, and moreover, by continuity,  $F_1$  and  $F_2$  are indistinguishable from these modifications.

We turn to the proof of (A.1). For an arbitrary fixed  $\gamma \in \Gamma$ , define the time-changed ‘bridge’ process  $X_\gamma$  by

$$X_\gamma(u) := M_\gamma(u) - \frac{V_\gamma(u)}{V_\gamma(1)} M_\gamma(1), \quad u \in [0, 1].$$

Then  $X_\gamma$  and  $M_\gamma(1)$  are independent, for they are jointly Gaussian with covariance function

$$\text{Cov}(X_\gamma(u), M_\gamma(1)) = V_\gamma(u) - \frac{V_\gamma(u)}{V_\gamma(1)} V_\gamma(1) = 0, \quad u \in [0, 1].$$

In terms of  $X_\gamma$  and  $M_\gamma(1)$ , we find

$$\tau_{\gamma 1} = \frac{1}{2} \frac{M_\gamma(1)^2 - V_\gamma(1)}{\int_0^1 M_\gamma^2(u) du} = \frac{1}{2} \frac{M_\gamma(1)^2 - V_\gamma(1)}{\int_0^1 (X_\gamma(u) + M_\gamma(1)q_\gamma(u))^2 du}$$

and

$$\tau_{\gamma 2} = \frac{1}{2} \frac{M_\gamma(1)^2 - V_\gamma(1)}{\sqrt{V_\gamma(1) \int_0^1 (X_\gamma(u) + M_\gamma(1)q_\gamma(u))^2 du}},$$

for  $q_\gamma(u) := V_\gamma(u)/V_\gamma(1)$ . The equality

$$P(\tau_{\gamma i} = x) = E[P(\tau_{\gamma i} = x | M_\gamma(1))] = 0$$

will hold for  $i = 1, 2$  and any  $x \in \mathbb{R}$  iff

$$P(\tau_{\gamma i} = x | M_{\gamma}(1)) = 0 \text{ a.s.}$$

for  $i = 1, 2$  and any  $x \in \mathbb{R}$ . In its turn, using the independence of  $X_{\gamma}(u)$  and  $M_{\gamma}(1)$ , the latter will hold if

$$P\left(\frac{1}{2} \frac{\alpha^2 - V_{\gamma}(1)}{\int_0^1 (X_{\gamma}(1) + \alpha q_{\gamma}(u))^2 \mathbf{d}u} = x\right) = 0,$$

$$P\left(\frac{1}{2} \frac{\alpha^2 - V_{\gamma}(1)}{\sqrt{V_{\gamma}(1)} \int_0^1 (X_{\gamma}(u) + \alpha q_{\gamma}(u))^2 \mathbf{d}u} = x\right) = 0$$

hold for all  $x \in \mathbb{R}$  and  $\alpha \neq \pm \sqrt{V_{\gamma}(1)}$  (because  $P(M_{\gamma}^2(1) = V_{\gamma}(1)) = 0$ ), which in its turn will hold if

$$P\left(\int_0^1 (X_{\gamma}(u) + \alpha q_{\gamma}(u))^2 \mathbf{d}u = x\right) = 0$$

for any  $\alpha, x \in \mathbb{R}$ . Since  $X_{\gamma}$  is a zero-mean Gaussian process with a continuous covariance and  $q_{\gamma}$  is square integrable, the equality in the previous display indeed holds, by o (a, pp. 472–473).  $\square$

The second lemma in this section allows to combine the conditional convergence of a Gaussian bootstrap process with a marginal convergence on the space of the data into a conditional convergence of a pair.

**Lemma A.2.** *Let the data be  $D_n = (M_n, U_n)$  and let the bootstrap multipliers be  $W_n^* = (w_1^*, \dots, w_n^*)'$ , with  $D_n$  independent of  $W_n^*$ . Let  $(M_n^*, X_n)$  be random elements of  $\mathcal{D}[0, 1] \times \mathcal{S}$  for some complete and separable metric space  $\mathcal{S}$ , such that  $M_n^*$  and  $X_n$  are measurable respectively w.r.t.  $(D_n, W_n^*)$  and  $D_n$ . Assume that  $M_n^*$  is, conditionally on the data, a zero-mean Gaussian process with independent increments and conditional variance function*

$$V_n^* = \phi(D_n, G_n) + o_p(1),$$

whereas  $X_n = \psi(D_n) + o_p(1)$  for some continuous functions  $\phi : \mathcal{D}_3[0, 1] \rightarrow \mathcal{D}[0, 1]$ ,  $\psi : \mathcal{D}_2[0, 1] \rightarrow \mathcal{S}$  and for some  $G_n \in \mathcal{D}[0, 1]$  satisfying  $G_n \rightarrow G$  in  $\mathcal{D}[0, 1]$  for a continuous  $G \in \mathcal{D}[0, 1]$ . Then under Assumptions 1 and 2 it holds that

$$(M_n^*, X_n) \xrightarrow{w^*} (M^*, X) | (M, U),$$

where  $M^*$  conditionally on  $(M, U)$  is a zero-mean Gaussian process with independent increments and conditional variance function  $\phi(M, U, G)$ , whereas  $X = \psi(M, U)$ .

PROOF OF LEMMA A.2. It holds that  $(D_n, V_n^*) \xrightarrow{w} (M, U, \phi(M, U, G))$  in  $\mathcal{D}_3[0, 1]$  by Lemma 1 and the CMT. Let  $v_n$  be measurable functions such that  $V_n^* = v_n(D_n)$ . Based on Theorem A.1, consider a Skorokhod representation  $\tilde{D}_n \stackrel{d}{=} D_n$  and  $(\tilde{M}, \tilde{U}) \stackrel{d}{=} (M, U)$  such that  $(\tilde{D}_n, v_n(\tilde{D}_n)) \xrightarrow{a.s.} (\tilde{M}, \tilde{U}, \phi(\tilde{M}, \tilde{U}, G))$  in  $\mathcal{D}_3[0, 1]$ .

On the added factor space of a product extension of the Skorokhod representation space, define  $\tilde{W}_n^* \stackrel{d}{=} W_n^*$ ; then  $\tilde{W}_n^*$  is independent of  $\tilde{D}_n$ . If  $\mu_n$  are measurable functions such that  $M_n^* = \mu_n(D_n, W_n^*)$ , define  $\tilde{M}_n^* = \mu_n(\tilde{D}_n, \tilde{W}_n^*)$ . Conditionally on  $\tilde{D}_n$ , the process  $\tilde{M}_n^*$  is a zero-mean Gaussian process with independent increments and conditional variance function  $v_n(\tilde{D}_n)$ . This holds because the conditional distribution of  $\tilde{M}_n^*$ , and the functions  $v_n$  in particular, are determined by the distribution of  $(\tilde{D}_n, \tilde{W}_n^*) \stackrel{d}{=} (D_n, W_n^*)$ . By construction, the conditional variance function of  $\tilde{M}_n^*$  satisfies  $v_n(\tilde{D}_n) \xrightarrow{a.s.} \phi(\tilde{M}, \tilde{U}, \tilde{G})$ . By fixing the outcomes in an appropriate measure-one set in the factor space of  $\tilde{D}_n$ , it follows by an outcome-by-outcome argument that  $\tilde{M}_n^* \xrightarrow{w}_{a.s.} \tilde{M}^* | (\tilde{M}, \tilde{U})$ , where  $\tilde{M}^*$  conditionally on  $(\tilde{M}, \tilde{U})$  is a zero-mean Gaussian process with independent increments and conditional variance function  $\phi(\tilde{M}, \tilde{U}, \tilde{G})$ . The convergence facts  $\tilde{D}_n \xrightarrow{a.s.} (\tilde{M}, \tilde{U})$  and  $\tilde{M}_n^* \xrightarrow{w}_{a.s.} \tilde{M}^* | (\tilde{M}, \tilde{U})$  jointly imply, by Lemma A.3 of a (a), the convergence

$$(\tilde{M}_n^*, \tilde{D}_n) \xrightarrow{w^*}_p (\tilde{M}^*, \tilde{M}, \tilde{U}) | (\tilde{M}, \tilde{U}) \quad (\text{A.2})$$

on the Skorokhod representation space (in fact, by the proof of the aforementioned Lemma A.3, also  $\xrightarrow{w^*}_{a.s.}$ ).

Finally, if the measurable functions  $\xi_n$  are such that  $X_n = \xi_n(D_n)$ , then  $\xi_n(\tilde{D}_n) = \psi(\tilde{D}_n) + o_p(1)$  because this equality is determined by the joint distribution of  $(\tilde{D}_n, \tilde{X}_n) \stackrel{d}{=} (D_n, X_n)$ . As  $\psi$  is continuous and upon conditioning convergence in probability to zero becomes weak convergence in probability to zero, from (A.2) and Theorem 10 of e (w) it follows that

$$(\tilde{M}_n^*, \xi_n(\tilde{D}_n)) \xrightarrow{w^*}_p (\tilde{M}^*, \psi(\tilde{M}, \tilde{U})) | (\tilde{M}, \tilde{U}).$$

The distributional equalities  $(M_n^*, X_n, D_n) \stackrel{d}{=} (\tilde{M}_n^*, \xi_n(\tilde{D}_n), \tilde{D}_n)$  and  $(M^*, X, M, U) \stackrel{d}{=} (\tilde{M}^*, \psi(\tilde{M}, \tilde{U}), \tilde{M}, \tilde{U})$  complete the proof.  $\square$

## A.2 Proofs

PROOF OF LEMMA 1. We follow the approach of the proof of Lemma 1 and other intermediate results in a (a). First, defining  $e_t = z_t^2 - 1$ ,

$$\sup_{u \in [0, 1]} |U_n(u) - V_n(u)| = \sup_{u \in [0, 1]} \left| n^{-1} \sum_{t=1}^{\lfloor nu \rfloor} \sigma_t^2 e_t \right| \xrightarrow{p} 0$$

by Theorem A.1 of v (a), since  $\{e_t, \mathcal{F}_t\}_{t \geq 1}$  is an mds by Assumption 1 and  $\sigma_{\lfloor n \cdot \rfloor + 1}^2 = \sigma_n^2(\cdot) \xrightarrow{w} \sigma^2(\cdot)$  by Assumption 2 and the CMT; this proves (8), because convergence in the sup norm

implies convergence in the Skorokhod metric, i.e., in  $\mathcal{D}[0, 1]$ . Next, we apply Theorem 2.1 of (a) to

$$M_n(\cdot) = \int_0^\cdot \sigma_n(u) dB_{z,n}(u),$$

noting that Assumption 1 implies  $\sup_{n \geq 1} n^{-1} \sum_{t=1}^n E(z_t^2) = 1$ , so that using Assumption 2, we have

$$(\sigma_n(\cdot), B_{z,n}(\cdot), M_n(\cdot)) \xrightarrow{w} (\sigma(\cdot), B_z(\cdot), M(\cdot)).$$

The CMT together with (8) then implies (7), because

$$\int_0^u \sigma_n^2(s) ds = \frac{1}{n} \sum_{t=1}^{\lfloor nu \rfloor} \sigma_t^2 + \sigma_{\lfloor nu \rfloor + 1}^2 (u - \lfloor nu \rfloor n^{-1}), \quad u \in [0, 1],$$

so that  $U_n(\cdot) = V_n(\cdot) + o_p(1) = \int_0^\cdot \sigma_n^2(s) ds + o_p(1)$ , i.e.,  $U_n(\cdot)$  is a continuous functional of  $\sigma_n(\cdot)$  plus an asymptotically negligible term.  $\square$

**PROOF OF THEOREM 1.** The idea of the proof is to construct on a special probability space random elements distributed like  $(\sigma_n, M_n, U_n, M_n^*, U_n^*)$  and such that on this probability space the convergence asserted in Theorem 1 holds weakly a.s.; on a general probability space it will then hold  $\xrightarrow{w}$ . Throughout, we use repeatedly the fact that for independent random elements  $\xi$  and  $\eta$  and for a measurable real  $\phi$  such that  $E(|\phi(\xi, \eta)|) < \infty$ , it holds that  $E(\phi(\xi, \eta)|\eta) = E(\phi(\xi, v))|_{v=\eta}$  a.s., with  $E(\phi(\xi, v))$  defining a function of a non-random  $v$ ; see [1, p. 341].

By Assumption 3,  $\psi_{nt}$  are  $\mathcal{G}_{n0}$ -measurable and hence are measurable functions of  $\sigma_n$  that we denote, with a slight abuse of notation, by  $\psi_{nt}(\sigma_n)$ . Let

$$e_{nm}(\gamma) := E \left( v_{nt}^2 \psi_{nt}^2(\gamma) \mathbb{I}_{\{|v_{nt} \psi_{nt}(\gamma)| > \sqrt{n}/m\}} \right),$$

for  $m \in \mathbb{N}$  and a generic non-random  $\gamma$ ; then  $e_{nm}(\sigma_n)$  is a version of the conditional expectation  $E \left( z_t^2 \mathbb{I}_{\{|z_t| > \sqrt{n}/m\}} | \sigma_n \right)$  because  $\{v_{nt}\}_{t=1}^n$  and  $\sigma_n$  are independent. Define  $B_{v,n} := n^{-1/2} \sum_{t=1}^{\lfloor n \cdot \rfloor} v_{nt}$ . We apply Theorem A.1 with  $\xi_n = (\sigma_n, B_{v,n})$ ,  $\xi = (\sigma, B_z)$ ,

$$f_n(\xi_n) = (\sigma_n, Q_{\psi,n}, Q_{z,n}, \mathcal{L}_n, L_n) \text{ and } f(\xi) = (\sigma, Q, Q, 0^\infty, 0^\infty),$$

where  $Q_{\psi,n} = n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \psi_{nt}^2$ ,  $Q_{z,n} = n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} z_t^2$ ,  $\mathcal{L}_n = \{n^{-1} \sum_{t=1}^n e_{nm}(\sigma_n)\}_{m \in \mathbb{N}} \in \mathbb{R}^\infty$ ,  $L_n = \{n^{-1} \sum_{t=1}^n z_t^2 \mathbb{I}_{\{|z_t| > \sqrt{n}/m\}}\}_{m \in \mathbb{N}} \in \mathbb{R}^\infty$ ,  $Q(u) = u, u \in [0, 1]$ , and  $0^\infty$  is the zero sequence in  $\mathbb{R}^\infty$ , the Frechet space. The functions  $f_n$  and  $f$  are defined on subspaces of the Borel space  $\mathcal{D}_2[0, 1]$  with the Skorokhod metric and the induced Borel  $\sigma$ -algebra, and take values in the Polish space  $\mathcal{D}_3[0, 1] \times \mathbb{R}^\infty \times \mathbb{R}^\infty$  with the product of the Skorokhod and the Frechet metric. The assumptions of the lemma imply  $(Q_{\psi,n}, Q_{z,n}) \xrightarrow{p} (Q, Q)$ , because  $(Q_{\psi,n} - Q, Q_{z,n} - Q)$  is the partial sum process of  $n^{-1}(\psi_{nt}^2 - 1, z_t^2 - 1)$ , which is an mda with respect to  $\mathcal{F}_t$  since

$E(\psi_{nt}^2 | \mathcal{F}_{t-1}) = E(z_t^2 | \mathcal{F}_{t-1}) = 1$  by the tower property; this partial sum converges to the zero function in probability by the corollary to Theorem 3.3 of n (a). Noting that, by applying Markov's conditional inequality,  $L_n \xrightarrow{P} 0^\infty$  follows from the corresponding result for  $\mathcal{L}_n = E(L_n | \mathcal{G}_{n0})$ , the assumptions of the lemma eventually imply  $f_n(\xi_n) \xrightarrow{w} f(\xi)$ .

Theorem A.1 then implies the existence of  $\tilde{\xi}_n = (\tilde{\sigma}_n, \tilde{B}_{v,n}) \stackrel{d}{=} (\sigma_n, B_{v,n})$  and  $\tilde{\xi} = (\tilde{\sigma}, \tilde{B}_z) \stackrel{d}{=} (\sigma, B_z)$ , defined on a single probability space and such that

$$\left( \tilde{\sigma}_n, \tilde{Q}_{\psi,n}, \tilde{Q}_{z,n}, \tilde{\mathcal{L}}_n, \tilde{L}_n \right) := f_n(\tilde{\xi}_n) \xrightarrow{a.s.} f(\tilde{\xi}) = (\tilde{\sigma}, Q, Q, 0^\infty, 0^\infty). \quad (\text{A.3})$$

Finally, we complete the set up by introducing a product extension of the previous probability space with generic outcomes  $(\tilde{\omega}, \omega^*)$  where a sequence  $\{\tilde{w}_t^*(\omega^*)\} \stackrel{d}{=} \{w_t^*\}$  and a standard Brownian motion  $\tilde{B}_z^*(\omega^*)$  are defined; these are thus independent of  $\{(\tilde{\sigma}_n, \tilde{B}_{v,n})\}_{n \geq 1}$  and  $(\tilde{\sigma}, \tilde{B}_z)$ .

As  $\tilde{B}_{v,n}$  and  $\tilde{\sigma}_n$  are independent (because  $B_{v,n}$  and  $\sigma_n$  are), it holds for any integrable random variable  $h(\tilde{\sigma}_n, \tilde{B}_{v,n})$  that  $E(h(\tilde{\sigma}_n, \tilde{B}_{v,n}) | \tilde{\sigma}_n) = E(h(\gamma, \tilde{B}_{v,n}) | \gamma = \tilde{\sigma}_n)$ . A similar equality holds for the independent  $\tilde{B}_z$  and  $\tilde{\sigma}$ . Therefore, to prove any convergence of the form

$$E\left(h_n(\tilde{\sigma}_n, \tilde{B}_{v,n}) | \tilde{\sigma}_n\right) \xrightarrow{a.s.} E\left(h(\tilde{\sigma}, \tilde{B}_z) | \sigma\right), \quad (\text{A.4})$$

it is sufficient to prove that  $E(h_n(\gamma_n, \tilde{B}_{v,n})) \rightarrow E(h(\gamma, \tilde{B}_z))$  for all deterministic sequences  $\{\gamma_n\}_{n \geq 1}$  in some set  $\Gamma \subset \mathcal{D}_\infty[0, 1]$  such that  $P(\{\tilde{\sigma}_n\}_{n \geq 1} \in \Gamma) = 1$ . We now choose and fix  $\Gamma$ . Consider the outcomes  $\tilde{\omega}$  such that convergence (A.3) holds at  $\tilde{\omega}$  and, moreover,  $(\int_0^\cdot \gamma d\tilde{B}_z^*)|_{\gamma = \tilde{\sigma}(\tilde{\omega})} = (\int_0^\cdot \tilde{\sigma} d\tilde{B}_z^*)(\tilde{\omega}, \omega^*)$  up to indistinguishability w.r.t. the measure of  $\tilde{B}_z^*$ ; here  $\int_0^\cdot \gamma d\tilde{B}_z^*$  is a Wiener integral defined on the factor space of  $\tilde{B}_z^*$  with square-integrable  $\gamma \in \mathcal{D}[0, 1]$ , whereas  $\int_0^\cdot \tilde{\sigma} d\tilde{B}_z^*$  is an Itô integral defined on the product space. A measure-one set of such outcomes  $\tilde{\omega}$  exists; see e.g. Lemma 3.2 of k (a). Define  $\Gamma \subset \mathcal{D}_\infty[0, 1]$  as the set of sequences  $\{\tilde{\sigma}_n(\tilde{\omega})\}_{n \geq 1}$  corresponding to  $\tilde{\omega}$  in such a set, then  $P(\{\tilde{\sigma}_n\}_{n \geq 1} \in \Gamma) = 1$  as required.

As noted in Remark 4.4, we may recover  $(M_n, U_n)$  (and hence the original data  $D_n$ ) from  $(\sigma_n, B_{v,n})$  as some measurable transformation, say  $m_n(\sigma_n, B_{v,n})$ . Define accordingly  $(\tilde{M}_n, \tilde{U}_n) := m_n(\tilde{\sigma}_n, \tilde{B}_{v,n})$  (and analogously  $\tilde{D}_n$ ). With  $\tilde{z}_{nt} := \tilde{\psi}_{nt} \tilde{v}_{nt}$ , where  $\tilde{\psi}_{nt} = \psi_{nt}(\tilde{\sigma}_n)$  and

$$\tilde{v}_{nt} := n^{1/2} \left( \tilde{B}_{v,n}(t/n) - \tilde{B}_{v,n}((t-1)/n) \right),$$

define also the process  $\tilde{B}_{z,n} := n^{-1/2} \sum_{t=1}^{\lfloor n \cdot \rfloor} \tilde{z}_{nt} =: m_{z,n}(\tilde{\sigma}_n, \tilde{B}_{v,n})$ , such that

$$(\tilde{\sigma}_n, \tilde{B}_{z,n}, \tilde{M}_n, \tilde{U}_n) \stackrel{d}{=} (\sigma_n, B_{z,n}, M_n, U_n).$$

We proceed to the convergence of  $(\tilde{M}_n, \tilde{U}_n)$  conditional on  $\tilde{\sigma}_n$  and prove that

$$E\left(g(\tilde{B}_{z,n}, \tilde{M}_n, \tilde{U}_n) \middle| \tilde{\sigma}_n\right) \xrightarrow{a.s.} E\left(g(\tilde{B}_z, \tilde{M}, \tilde{V}) \middle| \tilde{\sigma}\right) \quad (\text{A.5})$$

for continuous bounded real  $g$  of matching domain; this convergence is of the form (A.4) with  $h_n = g \circ (m_{z,n}, m_n)$ . In so doing, for any random element  $Z = \phi(\tilde{\sigma}_n, \tilde{B}_{v,n})$  we write  $Z(\gamma_n)$  for  $\phi(\gamma_n, \tilde{B}_{v,n})$ ; e.g.,  $\tilde{B}_{z,n}(\gamma_n) = m_{z,n}(\gamma_n, \tilde{B}_{v,n})$ . By the discussion in the previous paragraph, (A.5) will follow from the standard weak convergence of  $(\tilde{B}_{z,n}(\gamma_n), \tilde{M}_n(\gamma_n), \tilde{U}_n(\gamma_n))$ , for all  $\{\gamma_n\}_{n \geq 1} \in \Gamma$ , that we establish next.

For  $\{\tilde{\sigma}_n\}_{n \in \mathbb{N}}$  replaced by a fixed  $\{\gamma_n\}_{n \geq 1} \in \Gamma$ ,  $\tilde{z}_{nt}(\gamma_n) = \psi_{nt}(\gamma_n)\tilde{v}_{nt}$  is an mda satisfying the conditions of w (r)'s functional central limit theorem. First,  $E(\psi_{nt}(\gamma_n)\tilde{v}_{nt}|\{\tilde{v}_{ni}\}_{i=1}^{t-1}) = \psi_{nt}(\gamma_n)E(\tilde{v}_{nt}|\{\tilde{v}_{ni}\}_{i=1}^{t-1}) = 0$  because the mda property of  $\tilde{v}_{nt}$  is inherited from the original probability space as  $\{\tilde{v}_{ni}\}_{i=1}^n \stackrel{d}{=} \{v_{ni}\}_{i=1}^n$ . Second,  $n^{-1} \sum_{t=1}^{\lfloor n \rfloor} E(\psi_{nt}^2(\gamma_n)\tilde{v}_{nt}^2|\{\tilde{v}_{ni}\}_{i=1}^{t-1}) = n^{-1} \sum_{t=1}^{\lfloor n \rfloor} \psi_{nt}^2(\gamma_n) = \tilde{Q}_{\psi,n}(\gamma_n) \rightarrow Q$ , where the first equality is again inherited from the original probability space, and the convergence by the definition of  $\Gamma$ . Third, as  $\tilde{\mathcal{L}}_n(\gamma_n) \rightarrow 0^\infty$  again by the choice of  $\Gamma$ , it holds that  $n^{-1} \sum_{t=1}^n e_{nm}(\gamma_n) \rightarrow 0$  for all  $m \in \mathbb{N}$ , which is equivalent to

$$n^{-1} \sum_{t=1}^n E\left(\tilde{z}_{nt}^2(\gamma_n) \mathbb{I}_{\{|\tilde{z}_{nt}(\gamma_n)| > \sqrt{n}/m\}}\right) \rightarrow 0, \quad m \in \mathbb{N},$$

by the definition of  $e_{nm}$  and implies the Lindeberg condition in its usual form

$$n^{-1} \sum_{t=1}^n E\left(\tilde{z}_{nt}^2(\gamma_n) \mathbb{I}_{\{|\tilde{z}_{nt}(\gamma_n)| > \sqrt{n}\epsilon\}}\right) \rightarrow 0$$

for all  $\epsilon > 0$ . Therefore,

$$\tilde{B}_{z,n}(\gamma_n) \xrightarrow{w} \tilde{B}_z^*,$$

in the sense that  $E(g(\tilde{B}_{z,n}(\gamma_n))) \rightarrow E(g(\tilde{B}_z^*))$  for continuous bounded real  $g$  with matching domain. For the same fixed  $\gamma_n$ , this in turn implies that

$$\tilde{M}_n(\gamma_n) = \int_0^\cdot \gamma_n(u) d\tilde{B}_{z,n}(u, \gamma_n) \xrightarrow{w} \int_0^\cdot \gamma(u) d\tilde{B}_z^*(u),$$

where  $\gamma = \lim \gamma_n$  exists in  $\mathcal{D}[0, 1]$  by the choice of  $\gamma_n$ . More precisely, by Theorem 2.1 of s (a), as  $\sup_{n \geq 1} \sum_{t=1}^n E(\tilde{z}_{nt}^2(\gamma_n)) = \sup_{n \geq 1} \tilde{Q}_{\psi,n}(1, \gamma_n) < \infty$ , the previous convergence holds jointly with that of  $\tilde{B}_{z,n}$ , such that  $E(g(\tilde{B}_{z,n}(\gamma_n), \tilde{M}_n(\gamma_n))) \rightarrow E(g(\tilde{B}_z^*, \int_0^\cdot \gamma d\tilde{B}_z^*))$  for continuous bounded real  $g$ . Furthermore, using

$$\begin{aligned} \tilde{U}_n &= n^{-1} \sum_{t=1}^{\lfloor n \rfloor} \tilde{\sigma}_t^2 \tilde{\psi}_{nt}^2 + n^{-1} \sum_{t=1}^{\lfloor n \rfloor} \tilde{\sigma}_t^2 \left( \tilde{z}_{nt}^2 - \tilde{\psi}_{nt}^2 \right) \\ &= \int_0^\cdot \tilde{\sigma}_n^2(u) d\tilde{Q}_{\psi,n}(u) + n^{-1} \sum_{t=1}^{\lfloor n \rfloor} \tilde{\sigma}_t^2 \left( \tilde{z}_{nt}^2 - \tilde{\psi}_{nt}^2 \right) + o(1) \end{aligned}$$

uniformly, it follows that  $\tilde{U}_n(\gamma_n) \xrightarrow{P} \int_0^\cdot \gamma^2(u) du$  by Theorem A.1 of v (a), since  $\tilde{z}_{nt}^2(\gamma_n) - \tilde{\psi}_{nt}^2(\gamma_n)$  is an mda. As convergence in probability to a constant is joint with any weak convergence of

random elements defined on the same probability space, the convergence

$$E \left[ g(\tilde{B}_{z,n}(\gamma_n), \tilde{M}_n(\gamma_n), \tilde{U}_n(\gamma_n)) \right] \rightarrow E \left[ g \left( \tilde{B}_z^*, \int_0^\cdot \gamma d\tilde{B}_z^*, \int_0^\cdot \gamma^2 \right) \right]$$

is true for continuous bounded real  $g$  and  $\{\gamma_n\}_{n \geq 1} \in \Gamma$ , with  $\lim_{n \rightarrow \infty} \gamma_n = \gamma$ . Recall that, by the choice of  $\Gamma$ , for  $\tilde{\omega}$  in a set of probability one it holds that  $\{\tilde{\sigma}_n(\tilde{\omega})\}_{n \geq 1} \in \Gamma$ ,  $\tilde{\sigma}_n(\tilde{\omega}) \rightarrow \tilde{\sigma}(\tilde{\omega})$  and

$$\left( \tilde{B}_z^*(\omega^*), \left( \int_0^\cdot \gamma d\tilde{B}_z^* \right) (\omega^*), \int_0^\cdot \gamma^2 \right) \Big|_{\gamma = \tilde{\sigma}(\tilde{\omega})} = \left( \tilde{B}_z^*(\omega^*), \left( \int_0^\cdot \tilde{\sigma} d\tilde{B}_z^* \right) (\tilde{\omega}, \omega^*), \int_0^\cdot \tilde{\sigma}^2(\tilde{\omega}) \right)$$

up to  $\tilde{B}_z^*$ -indistinguishability. Since  $\tilde{B}_z^*$  is independent of  $\tilde{\sigma}$ , the two previous displays jointly imply

$$E \left[ g(\tilde{B}_{z,n}, \tilde{M}_n, \tilde{U}_n) \Big| \tilde{\sigma}_n \right] \xrightarrow{a.s.} E \left[ g \left( \tilde{B}_z^*, \int_0^\cdot \tilde{\sigma} d\tilde{B}_z^*, \tilde{V} \right) \Big| \tilde{\sigma} \right].$$

The proof of (A.5) is completed by using the distributional equality  $(\tilde{B}_z, \tilde{M}, \tilde{V}) \stackrel{d}{=} (\tilde{B}_z^*, \int_0^\cdot \tilde{\sigma} d\tilde{B}_z^*, \tilde{V})$ .

We turn to the bootstrap processes. Define

$$\tilde{B}_{z,n}^* := n^{-1/2} \sum_{t=1}^{\lfloor n \cdot \rfloor} \tilde{z}_{nt} \tilde{w}_t^*, \quad \tilde{M}_n^* := n^{-1/2} \sum_{t=1}^{\lfloor n \cdot \rfloor} \tilde{\sigma}_t \tilde{z}_{nt} \tilde{w}_t^*, \quad \tilde{U}_n^* := n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \tilde{\sigma}_t^2 \tilde{z}_{nt}^2 \tilde{w}_t^{*2}.$$

Here we show that

$$E \left( g(\tilde{B}_{z,n}^*, \tilde{M}_n^*, \tilde{U}_n^*) \Big| \tilde{\sigma}_n, \tilde{B}_{v,n} \right) \xrightarrow{a.s.} E \left( g(\tilde{B}_z^*, \tilde{M}^*, \tilde{V}) \Big| \tilde{\sigma} \right)$$

for continuous bounded real  $g$ , where  $\tilde{B}_z^*$  is a standard Brownian motion independent of  $(\tilde{\sigma}, \tilde{B}_z)$ , and  $\tilde{M}^* := \int_0^\cdot \tilde{\sigma} d\tilde{B}_z^*$ . Given that  $\{\tilde{w}_t^*\}$  and  $(\tilde{\sigma}, \tilde{B}_z)$  are independent, as in the proof of (A.5), we could proceed by fixing  $\{(\gamma_n, b_n)\}_{n \geq 1} \in \Gamma B$ , where  $\Gamma B$  is an appropriate set with  $P((\tilde{\sigma}_n, \tilde{B}_{v,n})_{n \geq 1} \in \Gamma B) = 1$ , and then discuss the standard weak convergence of  $(\tilde{B}_{z,n}^*, \tilde{M}_n^*, \tilde{U}_n^*)$  as a transformation of  $(\gamma_n, b_n, \{\tilde{w}_t^*\})$  instead of  $(\tilde{\sigma}, \tilde{B}_z, \{\tilde{w}_t^*\})$ . Since now  $(\tilde{\sigma}_n, \tilde{B}_{v,n})$  and  $\{\tilde{w}_t^*\}$  are defined on a product space, we implement this equivalently by fixing outcomes  $\tilde{\omega}$  in the component space of  $(\tilde{\sigma}_n, \tilde{B}_{v,n})$  and letting the outcome in the component space of  $\{\tilde{w}_t^*\}$  be the only source of randomness. In what follows, fix an  $\tilde{\omega}$  in a probability-one set where convergence (A.3) holds.

Then

$$n^{-1/2} \sum_{t=1}^{\lfloor n \cdot \rfloor} \tilde{z}_{nt}(\tilde{\omega}) \tilde{w}_t^* \xrightarrow{w} B_z^*,$$

because  $n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} E[\tilde{z}_{nt}^2(\tilde{\omega})(\tilde{w}_t^*)^2] = n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \tilde{z}_{nt}^2(\tilde{\omega}) = Q_{z,n}(\tilde{\omega}) \rightarrow Q$  and

$$L_n(\tilde{\omega}) = \left\{ n^{-1} \sum_{t=1}^n \tilde{z}_{nt}^2(\tilde{\omega}) \mathbb{I}(|\tilde{z}_{nt}(\tilde{\omega})| > \sqrt{n}/m) \right\}_{m \in \mathbb{N}} \rightarrow 0^\infty$$

by the choice of  $\tilde{\omega}$ , such that the following Lindeberg condition holds for every  $m \in \mathbb{N}$ :

$$n^{-1} \sum_{t=1}^n E[\tilde{z}_{nt}^2(\tilde{\omega})(\tilde{w}_t^*)^2 \mathbb{I}(|\tilde{z}_{nt}(\tilde{\omega}) \tilde{w}_t^*| > \sqrt{n}/m)]$$



$$\begin{aligned}
&\leq n^{-1} \sum_{t=1}^n E[\tilde{z}_{nt}^2(\tilde{\omega})(\tilde{w}_t^*)^2 \mathbb{I}(|\tilde{z}_{nt}(\tilde{\omega})\tilde{w}_t^*| > \sqrt{n}/m, |\tilde{w}_t^*| \leq K)] \\
&\quad + n^{-1} \sum_{t=1}^n E[\tilde{z}_{nt}^2(\tilde{\omega})(\tilde{w}_t^*)^2 \mathbb{I}(|\tilde{z}_{nt}(\tilde{\omega})\tilde{w}_t^*| > \sqrt{n}/m, |\tilde{w}_t^*| > K)] \\
&\leq n^{-1} \sum_{t=1}^n \tilde{z}_{nt}^2(\tilde{\omega}) \mathbb{I}(|\tilde{z}_{nt}(\tilde{\omega})| > \sqrt{n}/(mK)) \\
&\quad + E[(\tilde{w}_1^*)^2 \mathbb{I}(|\tilde{w}_1^*| > K)] \cdot n^{-1} \sum_{t=1}^n \tilde{z}_{nt}^2(\tilde{\omega}) \\
&\xrightarrow[n \rightarrow \infty]{} E\{(\tilde{w}_1^*)^2 \mathbb{I}(|\tilde{w}_1^*| > K)\} \xrightarrow[K \rightarrow \infty]{} 0.
\end{aligned}$$

It follows that  $\tilde{M}_n^*(\tilde{\omega}) = n^{-1/2} \sum_{t=1}^{\lfloor n \cdot \rfloor} \tilde{\sigma}_t(\tilde{\omega}) \tilde{z}_{nt}(\tilde{\omega}) \tilde{w}_t^* \xrightarrow{w} \int_0^1 \tilde{\sigma}(\tilde{\omega}) d\tilde{B}_z^*$ . Further,

$$\begin{aligned}
\tilde{U}_n^*(\tilde{\omega}) &= n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \tilde{\sigma}_t^2(\tilde{\omega}) \tilde{z}_{nt}^2(\tilde{\omega}) \tilde{w}_t^{*2} \\
&= \tilde{U}_n(\tilde{\omega}) + n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \tilde{\sigma}_t^2(\tilde{\omega}) \tilde{z}_{nt}^2(\tilde{\omega}) (\tilde{w}_t^{*2} - 1) \xrightarrow{p} \tilde{V}(\tilde{\omega}),
\end{aligned}$$

using Theorem A.1 of v (a). Since  $\tilde{V}(\tilde{\omega})$  is non-random, the last two convergence facts are joint:

$$E \left[ g \left( \tilde{M}_n^*(\tilde{\omega}), \tilde{U}_n^*(\tilde{\omega}) \right) \right] \rightarrow E \left[ g \left( \tilde{M}^*(\tilde{\omega}), \tilde{V}(\tilde{\omega}) \right) \right]$$

for continuous and bounded real  $g$ . As in the first part of the proof, by the product structure of the probability space and since the set of considered outcomes  $\tilde{\omega}$  has probability one, the previous convergence implies that

$$E \left( g(\tilde{M}_n^*, \tilde{U}_n^*) | \tilde{\sigma}_n, \tilde{B}_{v,n} \right) \xrightarrow{a.s.} E \left( g(\tilde{M}^*, \tilde{V}) | \tilde{\sigma} \right),$$

and eventually, as  $(\tilde{M}^*, \tilde{V}, \tilde{\sigma}) \stackrel{d}{=} (\tilde{M}, \tilde{V}, \tilde{\sigma})$ , that

$$E \left( g(\tilde{M}_n^*, \tilde{U}_n^*) | \tilde{\sigma}_n, \tilde{B}_{v,n} \right) \xrightarrow{a.s.} E \left( g(\tilde{M}, \tilde{V}) | \tilde{\sigma} \right).$$

Notice that conditioning on  $(\tilde{\sigma}_n, \tilde{B}_{v,n})$  can be replaced by conditioning on  $\tilde{D}_n$  because  $(\tilde{M}_n^*, \tilde{U}_n^*)$  is a measurable function of  $(\tilde{\sigma}_n, \tilde{B}_{v,n})$  and  $\{\tilde{w}_t^*\}$ .

We can conclude from (A.5) and this result that

$$\left( E \left[ h(\tilde{M}_n, \tilde{U}_n) | \tilde{\sigma}_n \right], E \left[ g(\tilde{M}_n^*, \tilde{U}_n^*) | \tilde{D}_n \right] \right) \xrightarrow{a.s.} \left( E \left[ h(\tilde{M}, \tilde{V}) | \tilde{\sigma} \right], E \left[ g(\tilde{M}, \tilde{V}) | \tilde{\sigma} \right] \right)$$

for all continuous and bounded real  $h, g$ , whereas on a general probability space

$$(E[h(M_n, U_n) | \sigma_n], E[g(M_n^*, U_n^*) | D_n]) \xrightarrow{w} (E[h(M, V) | \sigma], E[g(M, V) | \sigma]), \quad (\text{A.6})$$

because  $(\tilde{\sigma}_n, \tilde{M}_n, \tilde{U}_n, \tilde{D}_n, \tilde{M}_n^*, \tilde{U}_n^*) \stackrel{d}{=} (\sigma_n, M_n, U_n, D_n, M_n^*, U_n^*)$ . This is precisely the definition of the joint  $\xrightarrow{w}$  convergence in the theorem.  $\square$

**PROOF OF COROLLARY 1.** From (A.6) with  $h = g = \tau$ , if the random cdf  $P(\tau(M, V) \leq \cdot | \sigma)$  a.s. has continuous sample paths, conditional validity of the bootstrap as in Corollary 1 follows from Corollary 3.2 of v (a).  $\square$

PROOF OF LEMMA 1. For any  $K \in \mathbb{R}$ , consider the continuous function  $g_K : \mathbb{R} \rightarrow [0, 1]$  defined by  $g_K(x) = \mathbb{I}_{(-\infty, K]}(x) + (K + 1 - x)\mathbb{I}_{(K, K+1]}$ . Then  $\mathbb{I}_{(-\infty, K]} \leq g_K \leq \mathbb{I}_{(-\infty, K+1]}$  and the convergence  $\tau_n^* \xrightarrow{w^*} \tau^* | \sigma$  implies that

$$F_n^*(K) \leq E^*(g_K(\tau_n^*)) \xrightarrow{w} E(g_K(\tau) | \sigma) \leq F^*(K + 1),$$

where  $F^*(K + 1) = P(\tau^* \leq K + 1 | \sigma)$ . Therefore, for all  $q \in (0, 1)$ ,

$$\liminf_{n \rightarrow \infty} P(F_n^*(K) \leq q) \geq P(F^*(K + 1) \leq q).$$

As a result,

$$\begin{aligned} \liminf_{n \rightarrow \infty} P(F_n^*(\tau_n) \leq q) &\geq \liminf_{n \rightarrow \infty} P(F_n^*(\tau_n) \leq q, \tau_n \leq K) \\ &\geq \liminf_{n \rightarrow \infty} P(F_n^*(K) \leq q, \tau_n \leq K) \\ &\geq \liminf_{n \rightarrow \infty} P(F_n^*(K) \leq q) - \lim_{n \rightarrow \infty} P(\tau_n > K) \\ &\geq P(F^*(K + 1) \leq q), \end{aligned}$$

since  $\tau_n \xrightarrow{p} -\infty$  means that  $\lim_{n \rightarrow \infty} P(\tau_n > K) = 0$  for all  $K \in \mathbb{R}$ . By Markov's inequality,

$$P(F^*(K + 1) \leq q) \geq 1 - q^{-1}E(F^*(K + 1)) = 1 - q^{-1}P(\tau^* \leq K + 1),$$

and the proof is completed by letting  $K \rightarrow -\infty$ . □

PROOF OF EQ. (23). Notice that

$$\begin{aligned} \hat{U}_n(\cdot) &= n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \left( \sum_{i=0}^{t-1} \psi_i \varepsilon_{t-i} \right)^2 \\ &= n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \sum_{i=0}^{t-1} \psi_i^2 \varepsilon_{t-i}^2 + 2n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \sum_{i=0}^{t-1} \sum_{j=0}^{i-1} \psi_i \psi_j \varepsilon_{t-i} \varepsilon_{t-j} \\ &=: a_{1n}(\cdot) + a_{2n}(\cdot), \end{aligned}$$

with  $a_{1n}(\cdot)$  and  $a_{2n}(\cdot)$  implicitly defined. First,  $a_{2n}(\cdot) = o_p(1)$  uniformly in  $\cdot \in [0, 1]$ , similarly to Lemma A.7 in a (a). Second,

$$a_{1n}(\cdot) = n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \varepsilon_t^2 \left( \sum_{i=0}^{\lfloor n \cdot \rfloor - t} \psi_i^2 \right) = \left( \sum_{i=0}^{\infty} \psi_i^2 \right) U_n(\cdot) + b_n(\cdot),$$

with

$$b_n(\cdot) := n^{-1} \sum_{t=1}^{\lfloor n \cdot \rfloor} \varepsilon_t^2 \left( \sum_{i=\lfloor n \cdot \rfloor - t + 1}^{\infty} \psi_i^2 \right).$$

Since the  $\psi_i$ 's are exponentially decaying, there exist constants  $C$  and  $\rho \in (0, 1)$  such that  $\sum_{i=\lfloor n \cdot \rfloor - t + 1}^{\infty} \psi_i^2 \leq C\rho^{\lfloor n \cdot \rfloor - t + 1}$ . Using the facts that  $\max_{t=1, \dots, n} \sigma_t^2 = O_p(1)$  by Assumption 2 and  $E(z_t^2) = 1$  by Assumption 1, it holds that

$$\begin{aligned} \sup_{u \in [0, 1]} b_n(u) &\leq Cn^{-1} \sup_{u \in [0, 1]} \sum_{t=1}^{\lfloor nu \rfloor} \sigma_t^2 z_t^2 \rho^{\lfloor n \cdot \rfloor - t + 1} \\ &\leq C \left( \max_{t=1, \dots, n} \sigma_t^2 \right) \left( n^{-1} \max_{t=1, \dots, n} z_t^2 \right) \sup_{u \in [0, 1]} \left( \sum_{t=1}^{\lfloor n \cdot \rfloor} \rho^{\lfloor n \cdot \rfloor - t + 1} \right) \\ &= O_p(1) o_p(1) \sum_{t=1}^n \rho^t = o_p(1). \end{aligned}$$

Hence,  $\hat{U}_n(\cdot) = (\sum_{i=0}^{\infty} \psi_i^2) U_n(\cdot) + o_p(1)$ . □

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