TIME SERIES REGRESSION WITH A UNIT ROOT AND INFINITE-VARIANCE ERRORS

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In [4] Chan and Tran give the limit theory for the least-squares coefficient in a random walk with i.i.d. (identically and independently distributed) errors that are in the domain of attraction of a stable law. This paper discusses their results and provides generalizations to the case of I(1) processes with weakly dependent errors whose distributions are in the domain of attraction of a stable law. General unit root tests are also studied. It is shown that the semiparametric corrections suggested by the author in other work [22] for the finite-variance case continue to work when the errors have infinite variance. Surprisingly, no modifications to the formulas given in [22] are required. The limit laws are expressed in terms of ratios of quadratic functionals of a stable process rather than Brownian motion. The correction terms that eliminate nuisance parameters are zero in the limit and involve multiple stochastic integrals that may be written in terms of the quadratic variation of the limiting stable process. Some extensions of these results to models with drifts and time trends are also indicated.

1. INTRODUCTION

Suppose \{y_t\} is generated by

\[ y_t = \beta y_{t-1} + u_t; \quad t = 1, \ldots, n \quad (1) \]

with

\[ \beta = 1 \]

from an initialization at \( t = 0 \) in which \( y_0 \) is any random variable. Interest centers on the least-squares estimate

\[ \hat{\beta} = \left( \sum_{i=1}^{n} y_{i-1}^2 \right)^{-1} \left( \sum_{i=1}^{n} y_i y_{i-1} \right) \]

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of $\beta$ in (1). Chan and Tran [4] investigate the asymptotic behavior of $\hat{\beta}$ as $n \to \infty$ and show that for a certain family of i.i.d. errors $u_i$ with infinite variance the limit distribution of $n(\hat{\beta} - 1)$ can be characterized in terms of a functional of a Lévy process. They assume that $\gamma_0 = 0$, that $u_i \in D(\alpha)$, $u_i$ is in the domain of attraction of a stable law with index $\alpha$ ($0 < \alpha < 2$), and that the limit law satisfies the scaling condition
\begin{equation}
 n^{-1/\alpha} (u_1 + \cdots + u_n) \Rightarrow u_1,
\end{equation}
where "$\Rightarrow$" signifies equality in distribution (both here and elsewhere in the paper). Chan and Tran show that under these conditions
\begin{equation}
 n(\hat{\beta} - 1) = \left( \int_0^1 U^- dU \right)^{\alpha} \overset{D}{\to} \left( \int_0^1 U^2 \right)^{\alpha},
\end{equation}
where $U^-$ denotes the left limit of $U$, $U(r)$ is a Lévy process on the space of CADLAG functions $D[0,1]$ and "$\Rightarrow$" signifies weak convergence of measures. Their proof uses the weak convergence (from Resnick [27]) of the component processes
\begin{equation}
 \left( a^{-1} \sum_{i=1}^{[nr]} u_i, a^{-2} \sum_{i=1}^{[nr]} u_i^2 \right) \Rightarrow (U(r), V(r)),
\end{equation}
where $(U(r), V(r))$ is a Lévy process in $D[0,1]^2$ and the normalization is $a_n = n^{1/\alpha} \ell(n)$
\begin{equation}
\text{for some slowly varying function } \ell(n). \text{ They also show that}
\end{equation}
\begin{equation}
 \left( \frac{1}{2} \right) \left[ U^2(1) - V(1) \right] = \int_0^1 U^- dU
\end{equation}
and then (4) and (5) together with the continuous mapping theorem give the final result (3).

We shall start with some remarks on these results.

(i) Suppose $u_i$ belongs to what is known as the normal domain of attraction of a stable law with index $\alpha$ (see [12] Chapter 2 for a discussion of normal domains of attraction). We shall denote this by writing
\begin{equation}
 u_i \in \mathcal{D}(\alpha).
\end{equation}

The tail behavior of $u_i$ when $0 < \alpha < 2$ is then of the Pareto–Lévy form
\begin{equation}
 P(u_1 < u) = \frac{c_1 a^{\alpha}}{|u|^{\alpha}} \left[ 1 + o(1) \right], \quad u < 0
\end{equation}
\begin{equation}
 P(u_1 > u) = \frac{c_2 a^{\alpha}}{u^{\alpha}} \left[ 1 + o(1) \right], \quad u > 0
\end{equation}
as $|u| \to \infty$ ([12], p. 92). Here $c_1$ and $c_2$ are constants with $c_1, c_2 \geq 0$ and $c_1 + c_2 = 1$ (by suitable selection of $a$). We shall call $a$ the scale parameter.
In this case the norming sequence in (5) is of the simple form \( a_n = an^{1/\alpha} \) so that \( f(n) = a \) in (5).

(ii) Second, when (7) applies we have \( u_i^2 \in \mathcal{LD}(\alpha/2) \). Moreover, (4) can be replaced with an explicit limit law given as follows in terms of a stable process \( U_\alpha \):

\[
\left[ a_n^{-1} \sum_{i=1}^{[nr]} u_i, a_n^{-2} \sum_{i=1}^{[nr]} u_i^2 \right] = \left[ U_\alpha(r), \int_0^r (dU_\alpha)^2 \right],
\]

(10)

where \( a_n = an^{1/\alpha} \). Here \( U_\alpha(r) \) is a standard stable process with index \( \alpha \) and unit scale coefficient. When \( u_i = -u_i \) (so that the distribution of \( u_i \) is symmetric), \( U_\alpha(r) \) is a symmetric stable process and the characteristic function of \( U_\alpha(1) \) has the form \( e^{-ct|s|^\alpha} \) where

\[
c = \begin{cases} 
\Gamma(1 - \alpha) \cos(\pi\alpha/2), & \alpha \neq 1 \\
\pi/2, & \alpha = 1
\end{cases}
\]

[12, pp. 44–45]. Moreover, \( U_\alpha(r) = r^{1/\alpha} U_\alpha(1) \). These and other properties of stable processes are derived in Ito [13, pp. 157–162]. Figures 1–4 in the Appendix display typical trajectories of stable processes for various values of \( \alpha \). These may be contrasted with that of a typical Wiener process (Figure 5). We shall henceforth write

\( U_\alpha(r) = \text{SP}(\alpha) \)

to signify that \( U_\alpha \) is a standard stable process (SP) with index \( \alpha \). Note that the class \{SP(\alpha) : 0 < \alpha \leq 2 \} is a subclass of the Lévy processes and that each member of SP(\alpha) has no Wiener component in its Ito representation [27, p. 72] when \( \alpha < 2 \). (Thus, \( S = 0 \) and \( a = 0 \) in equation (10) of [4].)

In the representation (10) above, \( \int_0^r (dU_\alpha)^2 \) is a multiple stochastic integral which represents the usual quadratic variation (or square bracketed) process. This is sometimes represented in the notation \([U]_r\), (e.g. [19, p. 175]). But we prefer the integral notation in the present context because it helps to make the weak convergence that is given in (10) more intuitive and easily understood. In fact, Resnick proves (10) in [27, p. 94] but uses a notation for the limit in terms of point processes rather than stochastic integrals.

(iii) In place of the distributional equivalence of (6) [Theorem 2(ii) of [4]] we have indeed the direct equation

\[
V(1) = \int_0^1 (dU_\alpha)^2 = U_\alpha^2(1) - 2 \int_0^1 U_\alpha^- dU_\alpha.
\]

(11)

This follows from the Ito calculus for semimartingales (see, e.g., Kopp's second formula on page 160 of [15]).
It is most easily understood by noting that the stochastic differential $dU_\alpha^2$ can be broken down as follows:

$$dU_\alpha^2 = (U_\alpha + dU_\alpha)^2 - (U_\alpha)^2$$

$$= 2U_\alpha dU_\alpha + (dU_\alpha)^2.$$ 

Integration then yields formula (11) directly.

(iv) Finally, we note that when $\alpha = 2$,

$$U(r) = W(r) = \text{BM}(1)$$

or standard Brownian motion (BM). In this case

$$(dU_\alpha)^2 = (dW)^2 = dr$$

and (11) reduces to the usual formula

$$\int_0^1 W dW = \left( \frac{1}{2} \right) \left[ W^2(1) - 1 \right]$$

for the Brownian motion stochastic integral. This is, in fact, the only case for which $(dU_\alpha)^2$ is nonrandom. Note also that since $u_1 \in \mathcal{H}_D(\alpha)$ we necessarily have a finite variance $\sigma^2 = E(u_1^2) < \infty$ when $\alpha = 2$ (see [12], p. 92) and the scale factor is $a_n = on^{1/2}$. The limit distribution given in (3) is then the ratio of Brownian functionals

$$\int_0^1 W dW / \int_0^1 W^2.$$

2. GENERAL I(1) MODELS WITH INFINITE VARIANCE ERRORS

In econometrics there has recently been a good deal of interest in models such as (1) where allowance is made for some weak dependence in the errors $u_t$. The resulting time series are known as I(1) or integrated processes. My review paper [23] and Park and Phillips [20,21] provide a general discussion of models where these processes occur. In an earlier paper on scalar time series [22] I showed how to deal with such general error processes in constructing tests for the presence of a unit root. This involved a semiparametric correction to eliminate the bias in estimation of the regression coefficient that is due to the serial correlation in $u_t$. It is interesting to explore how this procedure needs to be modified when $u_t$ has infinite variance.

2.1. Models with MA(1) Errors

Let us start by considering the simple case of MA(1) errors

$$u_t = \varepsilon_t + \theta \varepsilon_{t-1}, \quad |\theta| < 1$$

(13)
where $\epsilon_i \in \mathcal{D}(\alpha)$, $0 < \alpha \leq 2$, and $\epsilon_i$ is i.i.d. with $\epsilon_i = -\epsilon_i$ (i.e., $\epsilon_i$ is symmetrically distributed). When $0 < \alpha < 2$ the tails of $\epsilon_i$ are of the Pareto-Lévy form with scale parameter $a$ and $c_1 = c_2 = \frac{1}{2}$. When $\alpha = 2$, $a^2 = E(\epsilon_i^2) < \infty$. Define $P_i = \sum_{t=1}^{i} \epsilon_t$ and then

$$
\sum_{1}^{n} y_{t-1} U_t = \sum_{1}^{n} P_{t-1} \epsilon_t + \theta^2 \sum_{1}^{n} P_{t-2} \epsilon_{t-1} + \theta \sum_{1}^{n} P_{t-1} \epsilon_{t-1} + \theta \sum_{1}^{n} P_{t-2} \epsilon_t
$$

$$
= \sum_{1}^{n} P_{t-1} \epsilon_t + \theta^2 \sum_{1}^{n} P_{t-2} \epsilon_{t-1} + \theta \sum_{1}^{n} \epsilon_{t-1}^2 + \theta \sum_{1}^{n} P_{t-2} (\epsilon_t + \epsilon_{t-1})
$$

$$
= \sum_{1}^{n} P_{t-1} \epsilon_t + \theta^2 \sum_{1}^{n} P_{t-2} \epsilon_{t-1} + \theta \left( \sum_{1}^{n} P_{t-1} \epsilon_t + \sum_{1}^{n} P_{t-2} \epsilon_{t-1} \right)
$$

$$
+ \theta \sum_{1}^{n} \epsilon_{t-1}^2 - \theta \sum_{1}^{n} \epsilon_{t-1} \epsilon_t.
$$

(14)

$$
\sum_{1}^{n} y_{t-1}^2 = \sum_{1}^{n} P_{t-1}^2 + \theta^2 \sum_{1}^{n} P_{t-2}^2 + 2\theta \sum_{1}^{n} P_{t-1} P_{t-2}.
$$

(15)

In the above we use the initialization $y_0 = 0$ to simplify formulae but this involves no loss of generality for the subsequent argument. Next observe that

$$
a^{-2} n^{-2/\alpha - 1} \sum_{1}^{n} P_{t-1}^2 = \int_{0}^{1} U_{\alpha}^2,
$$

(16)

$$
a^{-2} n^{-2/\alpha} \sum_{1}^{n} P_{t-1} \epsilon_t = \left( \frac{1}{2} \right) \left[ U_{\alpha}(1)^2 - \int_{0}^{1} (dU_{\alpha})^2 \right] = \int_{0}^{1} U_{\alpha}^- dU_{\alpha},
$$

(17)

and

$$
a^{-2} n^{-2/\alpha} \sum_{1}^{n} \epsilon_{t-1}^2 = \int_{0}^{1} (dU_{\alpha})^2,
$$

(18)

where $U_{\alpha}(r) = \text{SP}(\alpha)$, whereas

$$
a^{-2} n^{-2/\alpha} \sum_{1}^{n} \epsilon_{t-1} \epsilon_t \rightarrow 0.
$$

(19)

The latter follows because although the product $\epsilon_i \epsilon_{t-1} \in \mathcal{D}(\alpha/2)$ the cross product $\epsilon_i \epsilon_{t-1}$ does not lie in $\mathcal{D}(\alpha/2)$. In fact, $\epsilon_i \epsilon_{t-1} \in \mathcal{D}(\alpha)$ as shown by Cline [5] and Breiman [2] and (19) then follows directly. Note that the normalizing sequence for sums of these cross products is $b_n = b (n \ln n)^{1/\alpha}$ with $b = a^2$ as shown in Appendix A.

Standardizing (14) by $a^{-2} n^{-2/\alpha}$ and using (17)-(19) we have

$$
a^{-2} n^{-2/\alpha} \sum_{1}^{n} y_{t-1} U_t = (1 + \theta)^2 \int_{0}^{1} U_{\alpha}^- dU_{\alpha} + \theta \int_{0}^{1} (dU_{\alpha})^2.
$$

(20)
Similarly, using (16), (18) and (19) in (15) we get

\[ a^{-2} n^{-2/\alpha - 1} \sum_{i=1}^{n} y_{t-1}^2 = (1 + \theta)^2 \int_{0}^{1} U_{\alpha}^2. \]  \hspace{2cm} (21)

In view of (10), joint weak convergence of (20) and (21) also applies and we deduce the following limit result for the least squares estimator \( \hat{\beta} \)

\[ n(\hat{\beta} - 1) = \frac{(1 + \theta)^2 \int_{0}^{1} U_{\alpha}^- \, dU_{\alpha} + \theta \int_{0}^{1} (dU_{\alpha})^2}{(1 + \theta)^2 \int_{0}^{1} U_{\alpha}^2} \]

\[ = \frac{\int_{0}^{1} U_{\alpha}^- \, dU_{\alpha} + \theta (1 + \theta)^{-2} \int_{0}^{1} (dU_{\alpha})^2}{\int_{0}^{1} U_{\alpha}^2}. \]  \hspace{2cm} (22)

This formula generalizes (3) to the case of models with a unit root and MA(1) errors. Note that the second term in the numerator of (22) is random when \( \alpha < 2 \). When \( \alpha = 2 \) it is simply the constant \( \theta / (1 + \theta)^2 \). In that case (i.e. \( \alpha = 2 \)) the expression was given in my earlier paper [22, p. 283]. The effect of serially correlated errors in the unit root model (1) is therefore to induce a second order random bias term in the limit distribution of the least squares estimator. When the errors in (1) have finite variance this bias term is nonrandom and, as shown in [22], it depends on the serial correlation properties of the errors. The latter is still true in the infinite variance case but the bias term also has a random factor which depends on the quadratic variation \( \int_{0}^{1} (dU_{\alpha})^2 \).

The simplest way of dealing with the bias that is induced by MA(1) errors is to use instrumental variables (IV) estimation with \( y_{t-2} \) acting as an instrument for \( y_{t-1} \) in (1). Call the resulting estimator \( \tilde{\beta} = \sum_{2}^{n} y_{t} y_{t-2} / \sum_{2}^{n} y_{t-1} y_{t-2} \). It is easy to see that

\[ n(\tilde{\beta} - 1) = \int_{0}^{1} U_{\alpha}^- \, dU_{\alpha} / \int_{0}^{1} U_{\alpha}^2 \]

as in (3). IV estimators of this type have been suggested in the finite variance case by Hall [11] and Phillips and Hansen [26]. It is interesting to see that they continue to work as a direct method of eliminating the second order bias in the infinite variance case also.
2.2. Models with Weakly Dependent Errors

Let \( u_t \) be generated by the linear process

\[
\begin{align*}
  u_t &= d(L)\epsilon_t = \sum_{j=0}^{\infty} d_j \epsilon_{t-j}; \\
  d_0 &= 1, \\
  d(1) &\neq 0,
\end{align*}
\]

where \( \epsilon_t \) has the same properties as in (13). As shown in Brockwell and Davis [3, p. 480] (see also [14]) the series defining \( u_t \) converges almost surely if the coefficients \( d_j \) satisfy the condition

\[
\sum_{0}^{\infty} |d_j|^\delta < \infty, \quad \text{with } 0 < \delta < \alpha \wedge 1. \tag{24}
\]

If \( u_t \) is generated by a stable ARMA process then its moving average representation (23) has coefficients which decline geometrically, so that (24) is certainly satisfied in this case. It will be convenient for the limit theory below if we strengthen (24) to the following condition:

\[
\sum_{0}^{\infty} j |d_j|^\delta < \infty, \quad \text{with } 0 < \delta < \alpha \wedge 1. \tag{25}
\]

Again, this is satisfied by the coefficients of the moving average representation of a stable ARMA process.

THEOREM 2.1. If \( y_t \) is generated by (1) with \( \beta = 1 \), if \( u_t \) is the linear process (23) and if the summability condition (25) holds then

\[
\begin{align*}
  \left( \frac{1}{a_n} \sum_{1}^{n} u_t, a_n^{-2} \sum_{1}^{n} u_t^2 \right) &= \left( \omega U_n(1), \sigma^2 \int_{0}^{1} (dU_\alpha)^2 \right), \\
  a_n^{-2} \sum_{1}^{n} y_{t-1} u_t &= \omega^2 \int_{0}^{1} U_\alpha^{-1} dU_\alpha + \left( \frac{1}{2} \right) (\omega^2 - \sigma^2) \int_{0}^{1} (dU_\alpha)^2, \tag{27}
\end{align*}
\]

where \( a_n = an^{1/\alpha} \) and

\[
\begin{align*}
  \omega &= d(1) = \sum_{0}^{\infty} d_j, \\
  \sigma^2 &= \sum_{0}^{\infty} d_j^2. \tag{29}
\end{align*}
\]

Joint weak convergence of (27) and (28) also applies. \[ \blacksquare \]

The proof of Theorem 2.1 is given in Appendix B. The limit theory for the least squares estimator \( \hat{\beta} \) follows directly. We have:

\[
\begin{align*}
  n(\hat{\beta} - 1) &= \left( \int_{0}^{1} U_\alpha^2 \right)^{-1} \left( \int_{0}^{1} U_\alpha^{-1} dU_\alpha + \left( \frac{1}{2} \right) (1 - \sigma^2/\omega^2) \int_{0}^{1} (dU_\alpha)^2 \right) \tag{30}
\end{align*}
\]
generalizing both (3) and (22). Note that when $\alpha = 2$, (30) reduces to the expression derived earlier in [22, Theorem 3.1]. In that case the second term in the numerator of (30) is simply the constant $(\frac{1}{2})(\omega^2 - \sigma^2)$.

Similar results apply for other functionals like the $t$-ratio

$$t_{\beta} = (\hat{\beta} - 1)/s_{\beta}, \quad s_{\beta}^2 = n^{-1} \sum_1^n (y_x - \hat{\beta}y_{y-1})^2 \left( \sum_1^n y_x^2 \right)^{-1}. \quad (31)$$

Here we have

$$t_{\beta} = (\omega/\sigma) \left( \int_0^1 (dU_\omega)^2 \int_0^1 U_\omega^2 \right)^{-1/2}$$

$$\times \left( \int_0^1 U_\omega^- dU_\omega + \left( \frac{1}{2} \right) (1 - \sigma^2/\omega^2) \int_0^1 (dU_\omega)^2 \right). \quad (32)$$

Again this reduces to the formula given in [22, p. 282] for the finite variance case.

2.3. The Effect of Semiparametric Corrections

When $\{\epsilon_i\}$ is i.i.d., $(0,1)$ (i.e. $\alpha = 2$ and $\text{var}(\epsilon_i) = 1$) and $u_i$ is defined by (23) we see that

$$\sigma^2 = \text{var}(u_i), \quad \omega^2 = 2\pi f_u(0),$$

where $f_u(\cdot)$ is the spectral density of $u_i$. When $\alpha < 2$, the variance and the spectrum of $u_i$ are not finite because $E(\epsilon_i^2) = \infty$. The quantities $\sigma^2$ and $\omega^2$ do exist in this case at least as they are defined by (29). We shall call them pseudo-variances because they represent what would be the contribution to these variances in the usual formulae after the variance of $\epsilon_i$ is scaled out. These contributions remain finite even in the infinite variance case.

Similar remarks apply to the spectrum. When $\alpha = 2$ the autocorrelogram sequence for $u_i$ in (23) is given by

$$\rho(h) = E(u_iu_{i+h})/E(u_i^2) = \sum_1^\infty d_j d_{j+h}/\sum_1^\infty d_j^2, \quad h = 1,2,\ldots \quad (33)$$

Again the effects of $E(\epsilon_i^2)$ are scaled out and $\rho(h)$ is well defined as the final member of (33) when $\alpha < 2$ (see also [4] and [5] on this point). The $\rho$-spectrum of $u_i$ may then be defined as the Fourier transform of (33) i.e.

$$f^{(\rho)}_u(\lambda) = \frac{1}{2\pi} \sum_{h=-\infty}^\infty \rho(h)e^{i\lambda h}.$$
Davis and Resnick [7, 8] show that the sample autocorrelations are consistent for \( \rho(h) \) when \( \alpha < 2 \) i.e.

\[
 r(h) = \frac{\sum_{i=1}^{n-h} u_i u_{i+h}}{\sum_{i=1}^{n} u_i^2} - \rho(h)
\]

and that when \( \epsilon_t \) has Pareto–Lévy tails

\[
 r(h) - \rho(h) = O_p((\ln n/n)^{1/\alpha}).
\] (34)

In consequence, it is simple to construct consistent estimators of \( f_\nu^{(\rho)}(\lambda) \) using conventional spectral estimates based on the sample correlogram \( \{r(h): h = 1, 2, \ldots \} \).

In [22] I suggested some semiparametric corrections to \( n(\tilde{\beta} - 1) \) that asymptotically eliminate the nuisance parameter dependencies in the finite variance case. The statistic based on the coefficient estimate \( \tilde{\beta} \) has the form

\[
 Z(\tilde{\beta}) = n(\tilde{\beta} - 1) - \left( \frac{1}{2} \right) \left( \frac{n^{-2} \sum_{i=1}^{n} y_{i-1}^2}{n^{-1} \sum_{i=1}^{n} y_i^2 - \hat{\lambda}} \right)^{-1} (\hat{\omega}^2 - \hat{\sigma}^2),
\] (35)

where \( \hat{\omega}^2 \) and \( \hat{\sigma}^2 \) are consistent estimators of \( \omega^2 = 2\pi f_\nu(0) \) and \( \sigma^2 = E(u_i^2) \), respectively. Noting that in this case

\[
 \omega^2 = \sigma^2 + 2\lambda \text{ with } \lambda = \sum_{k=1}^{\infty} E(u_0 u_k)
\]

we may write \( Z(\tilde{\beta}) \) in the alternate form

\[
 Z(\tilde{\beta}) = \left( n^{-2} \sum_{i=1}^{n} y_{i-1}^2 \right)^{-1} \left( n^{-1} \sum_{i=1}^{n} y_i^2 - \hat{\lambda} \right),
\] (36)

where \( \hat{\lambda} \) is consistent for \( \lambda \). As shown in [22] in the finite variance model

\[
 Z(\tilde{\beta}) = \left( \int_0^1 W^2 \right)^{-1} \int_0^1 W dW
\] (37)

whose distribution is free of nuisance parameters. Thus, \( Z(\tilde{\beta}) \) forms the basis of a test for the presence of a unit root which is asymptotically similar for a wide class of weakly dependent errors with finite variance.

In the infinite variance (\( \alpha < 2 \)) case we have already seen that

\[
 a_n^{-2} \sum_{i=1}^{n} u_i^2 = \sigma^2 \int_0^1 (dU_\omega)^2, \text{ where } \sigma^2 = \sum_{i=0}^{\infty} d_i^2.
\]

Since \( \beta \) is consistent the same result holds if we replace \( u_t \) by the residuals \( \hat{u}_t = y_t - \tilde{\beta}y_{t-1} \). Thus we may write

\[
 na_n^{-2} \hat{\omega}^2 = \sigma^2 \int_0^1 (dU_\omega)^2.
\] (38)
Turning to $\omega^2$, we note that this parameter is usually estimated by a kernel procedure that leads to an expression of the general form

$$\hat{\omega}^2 = 2\pi f_n(0) = \sum_{j=-M}^{M} k(j/M)c(j),$$  

where

$$c(j) = n^{-1} \sum_{t=1}^{n} u_t u_{t+j}, \quad 1 \leq t + j \leq n$$

and the lag window $k(\cdot)$ is a bounded even function defined on the interval $[-1,1]$ with $k(0) = 1$. $M$ is a bandwidth parameter in (39) and it satisfies $M \to \infty$ and $M/n \to 0$ as $n \to \infty$. For example, when $k(j/M) = 1 - |j|/M$, $\hat{\omega}^2$ is the Bartlett estimator of what would be the long run variance of $u_t$, if $\alpha = 2$ (see [22] for further discussion).

Observe that for fixed $j$

$$na_n^{-2}c(j) = \left(\int_{0}^{1} dU_n\right)^2.$$

The same result also applies when $u_t$ is replaced by the residual $\hat{u}_t$ in (40). I shall not give a complete derivation here but using the same approach as that in the proof of Theorem 3.1 of my paper [24] it can be shown that, if $M = o(n^{1/2})$ as $n \to \infty$, then

$$na_n^{-2}\hat{\omega}^2 = \omega^2 \int_{0}^{1} (dU_n)^2.$$

Now note that we may write

$$Z(\hat{\beta}) = n(\hat{\beta} - 1) - \left(\frac{1}{2}\right) n^{-1} a_n^{-2} \left(\sum_{t=1}^{n} y_{t-1}^2\right)^{-1} \{na_n^{-2}(\hat{\omega}^2 - \hat{\delta}^2)\}.$$

From (31), (38) and (41) we deduce that

$$Z(\hat{\beta}) = \left(\int_{0}^{1} U_n^2\right)^{-1} \left(\int_{0}^{1} U_n^{-2} \, dU_n\right).$$

This result generalizes (37) to the infinite variance case and of course also includes (37) since $U_n(r) = W(r)$ when $\alpha = 2$.

The $t$-ratio statistic may be analyzed in the same way. The statistic I suggested in [22] is based on $t_\beta$ and has the form

$$Z(t) = (\hat{\delta}/\hat{\omega}) t_\beta - \left(\frac{1}{2}\right) \left(\hat{\omega}^2 - \hat{\delta}^2\right) \left\{\hat{\omega} \left(n^{-2} \sum_{t=1}^{n} y_{t-1}^2\right)^{1/2}\right\}^{-1}.$$
As \( n \to \infty \) we find that

\[
Z(t) = \left\{ \int_0^1 (dU_\alpha)^2 \int_0^1 U_\alpha^2 \right\}^{-1/2} \int_0^1 U_\alpha dU_\alpha.
\]  

(43)

When \( \alpha = 2 \) the limit distribution becomes

\[
\left( \int_0^1 W^2 \right)^{-1/2} \int_0^1 W dW
\]

as given in [22].

These results show that the semiparametric corrections suggested in [22] continue to work when the errors have infinite variance even though they were designed specifically to eliminate nuisance parameters in the finite variance case. The reason for this is that in the infinite variance case there are still parametric dependencies in the limit distributions of the coefficient estimator and its \( r \)-statistic (as shown in (30) and (32)). These dependencies involve the parameters \( \sigma^2 \), \( \omega^2 \) that we have described as pseudo-variances. They represent what would be the variance and the long-run variance of \( u_t \) if \( \alpha \) were equal to 2 and \( E(\varepsilon_t^2) < \infty \). The semiparametric corrections eliminate these pseudo-variances from the limit distributions in the infinite variance case just as they do when the actual variances are finite.

3. ADDITIONAL REMARKS

(i) The final results (42) and (43) apply under somewhat weaker assumptions than those given here. We may, for example, replace the requirement that \( \varepsilon_t \in \mathcal{D}(\alpha) \) with \( \varepsilon_t \in \mathcal{D}(\alpha) \). This affects the norming sequence \( a_n \) but we still find for a suitable choice of \( a_n \) that

\[
a_n^{-1} \sum_{t=1}^{[nr]} \varepsilon_t, a_n^{-2} \sum_{t=1}^{[nr]} \varepsilon_t^2 \Rightarrow \left( U_\alpha(r), \int_0^r (dU_\alpha)^2 \right)
\]

(44)

and

\[
a_n^{-2} \sum_{t=1}^n \varepsilon_t \varepsilon_{t+j} \Rightarrow 0; \quad 1 \leq t + j \leq n, \quad j \neq 0.
\]

These limits ensure that (38) and (41) hold, giving the final results (42) and (43) as stated.

(ii) We may also relax the symmetry condition \( \varepsilon_t = -\varepsilon_t \), although many of the arguments given in Section 2 will then need modification. When \( \alpha < 1 \) no further requirement beyond \( \varepsilon_t \in \mathcal{D}(\alpha) \) seems to be needed. When \( \alpha > 1 \) we require \( E(\varepsilon_t) = 0 \), as in [4], so that sums involving \( \varepsilon_t \) do not need to be centered. When \( \alpha = 1 \) an additional condition such as

\[
b_n = E(\varepsilon_t \mathbf{1}(|\varepsilon_t| \leq a_n)) = 0, \text{ for all } n
\]

(45)
ensures that centering of the sums involving $\epsilon_t$ is unnecessary, although this
is hardly weaker than the symmetry condition $\epsilon_t = -\epsilon_t$. An alternative proof
of Theorem 2.1 under such conditions will be reported elsewhere.

(iii) Our analysis and results extend easily to models with drifts and time
trends (or other deterministic functions) in place of (1). In such cases the time
series may be regarded as filtered prior to their use in regressions such as (1).
The effects can then be determined by treating the filtered series as regres-
sion residuals. For instance, when there is polynomial detrending we con-
struct the residual process $y_t$ from the least squares regression

$$y_t = \hat{\mu}_0 + \hat{\mu}_1 t + \ldots + \hat{\mu}_p t^p + \epsilon_t.$$  

Then in place of (27) and (28) we have

$$a_n^{-2} \sum_{i=1}^n y_{i-1} u_i = \omega^2 \int_0^1 U_a^{-1} dU_a + \left( \frac{1}{2} \right) (\omega^2 - \sigma^2) \int_0^1 (dU_a)^2 \quad (27)$$

and

$$n^{-1} a_n^{-2} \sum_{i=1}^n y_{i-1}^2 = \omega^2 \int_0^1 U_a^2 \quad (28)$$

where $U_a = QU_a$ is simply the projection of $U_a$ in $L_2[0,1]$ on the orthog-
ornal complement of the space spanned by the polynomial functions $\{0(r),
1(r), \ldots, p(r); j(r) = r^j\}$. We deduce that if $\hat{\beta}$ is the least squares estimator
of $\beta$ in

$$y_t = \beta_0 + \beta_1 t + \ldots + \beta_p t^p + \epsilon_t + u_t \quad (1)$$

under $\beta = 1$ and $\beta_p = 0$ then

$$n(\hat{\beta} - 1) = \left\{ \int_0^1 U_a^2 \right\}^{-1} \left\{ \int_0^1 U_a^{-1} dU_a + \left( \frac{1}{2} \right) (1 - \sigma^2/\omega^2) \int_0^1 (dU_a)^2 \right\}.$$

(30)

Semiparametric corrections to eliminate the nuisance parameters in (30) may
be made as in Section 2.3. The situation here is entirely analogous to that explo-
erd in the finite variance case by Park and Phillips [20, 21].

(iv) The limit theory given here also has applications in the context of limit
theorems for self normalized sums. These have been considered elsewhere
recently by several authors [1, 17, 25, 27]. In [25] the bimodality of $t$-ratio
statistics of the form

$$t_s = \frac{\sum_{i=1}^n \epsilon_i}{\left( \sum_{i=1}^n \epsilon_i^2 \right)^{1/2}}$$

was explored when $\epsilon_i \in \mathcal{M}_d(\alpha)$ and $0 < \alpha < 2$. The reason for
the bimodality in the distribution of $t_s$, which occurs in both finite and asymptotic sam-
samples, is the statistical dependence between the numerator and denominator random variables in $t_n$. When $e_j = -e_i$, $e_j \in \mathcal{D}(\alpha)$ and $e_j$ is i.i.d. it follows immediately from (44) that

$$t_n = U_\alpha(1) \left( \int_0^1 (dU_\alpha)^2 \right)^{1/2}$$

(46)

as $n \to \infty$. Using (11) we may now write the limit law in the form

$$U_\alpha(1) \left( \left( U_\alpha(1)^2 - 2 \int_0^1 U_\alpha^- dU_\alpha \right) \right)^{1/2}.$$  

(47)

The bimodality (with modes at $\pm 1$) then arises because of the occurrence of $U_\alpha(1)$ in the numerator and denominator elements of (47) when $0 < \alpha < 2$. When $\alpha = 2$ the denominator is nonrandom, of course, and (46) corresponds to the conventional limit theory for the $t$-ratio in the finite variance case, i.e. $t_n \Rightarrow U_2(1) = N(0,1)$.

REFERENCES


**APPENDIX A: ON THE TAIL BEHAVIOR OF THE PRODUCT $X = X_1X_2$ OF INDEPENDENT VARIATES $X_1 \in \mathcal{H}(\alpha)$**

Suppose $x_i = -x_i$ and then $X = -X$. Let $f_i(x)$ be the density of $x_i$, which we take to be continuous over $(-\infty, \infty)$. Setting $c_1 = c_2 = \frac{1}{2}$ in (8) and (9), we have the following tail behavior

$$f_i(x) = \left(\frac{1}{2}\right) \alpha a x |x|^{-\alpha - 1}(1 + o(1)), \quad |x| > k$$

for some (possibly large) constant $k > 0$. The density of $X$ is now

$$f(X) = 2 \int_0^{\infty} (1/x)f_1(x)f_2(X/x)dx$$

$$= 2 \left( \int_0^{\infty} + \int_k^{X/k} + \int_{X/k}^{\infty} \right) (1/x)f_1(x)f_2(X/x)dx.$$ 

Suppose $X > k^2$ and let $c$ be a generic constant in what follows. The first integral is less than

$$2c \int_0^{k} (1/x)f_2(X/x)dx = 2cX^{-1-\alpha} \left( \int_0^{k} c a x^{-\alpha}dx + o(1) \right) = O(X^{-1-\alpha}).$$
The third is less than
\[ 2c \int_{X/k}^\infty \frac{1}{x} f_1(x) \, dx = 2c a_1 a_2^\alpha \int_{X/k}^\infty x^{-\alpha-2} \, dx = O(X^{-1-\alpha}). \]

The second integral dominates and has the form
\[ \left( \frac{1}{2} \right) \alpha^2 a_1^2 X^{-1-\alpha} \int_k^{X/k} \frac{1}{x} \, dx = \left( \frac{1}{2} \right) \alpha^2 a_1^2 X^{-1-\alpha} \ln X (1 + o(1)) \]
so that \( f(X) = O(\ln X)X^{-1-\alpha} \) as \( |X| \to \infty \). Integrating again we get the tail behavior of the cumulative distribution function
\[ \text{cdf}(-X) = \left( \frac{1}{2} \right) f(\ln(-X)) (-X)^{-\alpha}(1 + o(1)) \]
\[ 1 - \text{cdf}(X) = \left( \frac{1}{2} \right) f(\ln X) X^{-\alpha}(1 + o(1)) \]
as \( X \to \infty \) where \( f = \alpha a_1^\alpha a_2^\beta \). Since \( \ln(\cdot) \) is slowly varying at infinity these tail probabilities ensure that \( X \in \mathcal{D}(\alpha) \) \cite[Theorem 2.6.1]{12}. Setting \( b_n = b(n \ln n)^{1/\alpha} \) with \( b = a_1 a_2 \) we have
\[ n[\text{cdf}(b_n X)] \to \left( \frac{1}{2} \right) (-X)^{-\alpha}, \quad X < 0 \]
\[ n[1 - \text{cdf}(b_n X)] \to \left( \frac{1}{2} \right) X^{-\alpha}, \quad X > 0. \]

The norming sequence for sums of i.i.d. variates distributed as \( X \) is therefore \( b_n = a_1 a_2 (n \ln n)^{1/\alpha} \) by the argument in \cite[p. 76]{12}. Explicitly we have \( b_n = \inf x |P(|X| > x) \leq 1/n| = a_1 a_2 (n \ln n)^{1/\alpha} \). Note that \( X \in \mathcal{D}(\alpha) \) and not \( \mathcal{S}(\alpha) \) because of the presence of the slowly varying function \( \ln(\cdot) \) in the tail formulae.

When the variates \( x_i \in \mathcal{S}(\alpha) \), \( \alpha_1 \neq \alpha_2 \) we obtain by a similar argument the result \( X \in \mathcal{S}(\alpha_1 \wedge \alpha_2) \). In this case the product variate \( X \) is in the normal domain of attraction of the law of the component variate with the smallest exponent.

**APPENDIX B: PROOF OF THEOREM 2.1**

Result (26) follows from Theorems 4.1 and 4.2 of Davis and Resnick \cite[pp. 189-192]{7}. Result (27) follows from (26) and summation by parts since
\[ a_n^{-2} \sum_{i=1}^n y_{i-1} u_i = \left( \frac{1}{2} \right) \left[ \left( a_n \sum_{i=1}^n u_i \right)^2 - a_n^{-2} \sum_{i=1}^n u_i^2 \right] \]
\[ = \left( \frac{1}{2} \right) \left[ \omega^2 U_n (1)^2 - \sigma^2 \int_0^1 (dU_n)^2 \right] \]
\[ = \omega^2 \int_0^1 U_n^- dU_n + \left( \frac{1}{2} \right) (\omega^2 - \sigma^2) \int_0^1 (dU_n)^2 \]
in view of (11). To prove (28) we use the decomposition

$$y_t = d(L)\epsilon_t = d(1)\epsilon_t + \tilde{\epsilon}_{t-1} - \tilde{\epsilon}_t,$$

where $\tilde{\epsilon}_t = d(L)\epsilon_t = \sum_{j=0}^{\infty} \tilde{d}_j \epsilon_{t-j}$ and $\tilde{d}_j = \sum_{j=1}^{\infty} d_k.$ Observe that

$$\sum_{j=0}^{\infty} |\tilde{d}_j|^b = \sum_{j=0}^{\infty} \left|\sum_{j=1}^{\infty} d_k\right|^b < \sum_{j=0}^{\infty} \sum_{k=j+1}^{\infty} |d_k|^b$$

$$= \sum_{k=j}^{\infty} |d_k|^b \sum_{j=0}^{k-1} 1$$

$$= \sum_{j=0}^{\infty} j |d_k|^b < \infty$$

under (25). Hence, $\tilde{\epsilon}_t = \sum_{j=0}^{\infty} \tilde{d}_j \epsilon_{t-j}$ converges almost surely and $\tilde{\epsilon}_t \in \mathcal{D}(\alpha).$ Thus, the decomposition (48) is well defined. Set $P_t = \sum_{i=1}^{t} \epsilon_i$ and then we have

$$n^{-1}a_n^{-2} \sum_{i=1}^{n} \epsilon_i^2 = \omega^2 \left(n^{-1}a_n^{-2} \sum_{i=1}^{n} P_{i-1}^2 + 2d(1) \left[ n^{-1}a_n^{-2} \sum_{i=1}^{n} P_{i-1} (\tilde{\epsilon}_{i-1} - \tilde{\epsilon}_i) \right] \right)$$

$$+ n^{-1}a_n^{-2} \sum_{i=1}^{n} (\tilde{\epsilon}_{i-1} - \tilde{\epsilon}_i)^2.$$

(49)

But

$$n^{-1}a_n^{-2} \sum_{i=1}^{n} P_{i-1}^2 = \int_{0}^{1} U_n^2$$

by (10) and the continuous mapping theorem. Thus, (28) follows if we can show that the final two terms on the right side of (49) tend in probability to zero. Consider

$$n^{-1}a_n^{-2} \sum_{i=1}^{n} P_{i-1} (\tilde{\epsilon}_{i-1} - \tilde{\epsilon}_i) = n^{-1}a_n^{-2} \left[ \sum_{i=1}^{n} P_{i-1} \tilde{\epsilon}_{i-1} - \sum_{i=1}^{n} (P_i - \epsilon_i) \tilde{\epsilon}_i \right]$$

$$= -n^{-1}a_n^{-2} \left[ P_n \tilde{\epsilon}_n - \sum_{i=1}^{n} \epsilon_i \tilde{\epsilon}_i \right]$$

$$= -n^{-1} \left[ (a_n^{-1} P_n)(a_n^{-1} \tilde{\epsilon}_n) - a_n^{-2} \sum_{i=1}^{n} \epsilon_i \tilde{\epsilon}_i \right]$$

$$= O_p(n^{-1})$$

since $\epsilon_i \tilde{\epsilon}_i \in \mathcal{D}(\alpha/2).$ Thus, the second term on the right of (49) converges in probability to zero. Finally, since $(\tilde{\epsilon}_t - \tilde{\epsilon}_{t-1})^2 \in \mathcal{D}(\alpha/2)$ we have

$$n^{-1}a_n^{-2} \sum_{i=1}^{n} (\tilde{\epsilon}_{i-1} - \tilde{\epsilon}_i)^2 = O_p(n^{-1})$$

and the third term on the right of (49) converges in probability to zero. \qed
Figure 1. Symmetric stable motion, $\alpha = 1.9$.

Figure 2. Symmetric stable motion, $\alpha = 1.5$. 
**Figure 3.** Symmetric stable motion, $\alpha = 1.0$.

**Figure 4.** Symmetric stable motion, $\alpha = 0.5$. 
Figure 5. Brownian motion ($\alpha = 2.0$).