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ASSESSING CONTEMPORANEOUS CORRELATION IN LONGITUDINAL DATA
by

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#### Abstract

Consider drawing a sample of ' $n$ ' experimental units where each unit is observed over ' $T$ ' time periods. Are the draws independent? Small and large (in ' $n$ ') sample properties of a distribution-free statistic, built to assess the degree of correlation, are provided. Further, the large sample distribution is also investigated in the case where the statistic is evaluated using residuals from a complex model.


## 1. INTRODUCTION

A strong case can be made for the statement that the regression model is the pre-eminent statistical model, at least in terms of the number of applications. Because arguments using regression models have become widely accepted in scientific communities, researchers are becoming more willing to collect regression data over time. If the data is identified at the level of the experimental unit then longitudinal, or sometimes called pooled time series cross-sectional, models are appropriate for fitting. At least in the social sciences, even though this data often consists of many experimental units, experimenters are often unwilling or unable to consider the data over long periods of time because basic conditions appear unstable.

In the life sciences, often the parent population is sufficiently large and randomization techniques for observation selection are sufficiently sound so that experimental units may be modeled independently. In the social sciences, however, there may exist important correlations between experimental units especially in applications where economic entities are the experimental unit. I follow standard terminology in econometrics and call these "contemporaneous" correlations to distinguish them from autocorrelations, or correlations through time. If detected, there are several modeling techniques available. Generalized least squares techniques and, in particular, seemingly unrelated regression models of the variance structure, are prominent in applications. Further, contemporaneous correlations can be induced by omitted covariates that are possibly common to the experimental units. A subtler cause of contemporaneous correlation can arise from the random nature of one or several parameters in the model of the mean effect. The point here is that there are several type of models to handle contemporaneous correlation when detected. This article concerns "omnibus" methods for detecting contemporaneous correlations, particularly when the number of experimental units, $n$, is large compared to the number of observations per experimental unit, T. Hsiao (1986) provides a detailed discussion of the importance of considering the case of large $n$ in addition to the usual approach in econometrics for considering large $T$.

Specifically, initially consider the model

$$
\begin{equation*}
Y_{i t}=\mu_{i}+\sigma_{i} e_{i t} \quad i=1, \ldots, n \quad t=1, \ldots, T \tag{1.1}
\end{equation*}
$$

where $Y_{i t}, t=1, \ldots, T$ are observations from the $i^{\text {th }}$ experimental unit. Assume $\left\{\mathrm{e}_{\mathrm{it}}\right\}$ are mean zero, unobservable random variables. Only the case of an equal number of observations per unit is explicitly discussed here as this is the rule when considering economic data (often violated, as with most rules). Many of the results presented in this article could be generalized to the unequal case with considerably more tedious algebra and at the expense of the interpretation of results. Both $\mu$ and $\sigma$ may be considered to be either parameters to be estimated, covariate effects, or some combination. It is interesting to note that it tums out that $\mu_{\mathrm{i}}$ may be either a fixed or random effect. It is well-known that location and scale parameters do not affect correlations and this also turns out to be
the case here. Model extensions are considered in Section 4.

Classical statistics for assessing the independence of experimental observations can be found in multivariate analysis. Perhaps the most well-known is Bartlett's (1954) statistic, defined as

$$
\begin{equation*}
C_{B}=\operatorname{det}(B)^{T / 2} / \prod_{i=1}^{\mathrm{n}} b_{i i}^{T / 2} \tag{1.2}
\end{equation*}
$$

where $B$ is a $n \times n$ matrix whose $(i, j)^{\text {th }}$ element is $b_{i j}=(T-1)^{-1} \Sigma_{\mathrm{t}}\left(\mathrm{Y}_{\mathrm{it}}-\overline{\mathrm{Y}}_{\mathrm{i}}\right)\left(\mathrm{Y}_{\mathrm{jt}}-\overline{\mathrm{Y}}_{\mathrm{j}}\right)$. This statistic is the maximum likelihood statistic under normality. It is straightforward to check that rank $(B)=$ minimum of ( $\mathrm{n}, \mathrm{T}$ 1) and thus $\operatorname{det}(B)=0$ for large $n$. The diagnostic statistics used in this paper are versions of a Lagrange Multiplier test statistic due to Breusch and Pagan (1980),

$$
\begin{equation*}
C_{A V E}^{2}=\left(\frac{n}{2}\right)^{-1} \Sigma_{i<j} c_{i j}^{2} \tag{1.3}
\end{equation*}
$$

where $c_{i j}=b_{i j} /\left(b_{i i} b_{j j}\right)^{1 / 2}$ is the Pearson correlation coefficient between the $i^{\text {th }}$ and $j^{\text {th }}$ flows. Breusch and Pagan (1980) showed that as $\mathrm{T} \rightarrow \infty$, the limiting distribution of (T-1) ( $\frac{\mathrm{n}}{2}$ ) $\mathrm{C}_{\mathrm{AVE}}^{2}$, is $\chi^{2}$ with $\mathrm{n}(\mathrm{n}-1) / 2$ degrees of freedom. This suggests the modified version $\mathrm{C}_{\mathrm{STD}}=\mathrm{n}\left((\mathrm{T}-1) \mathrm{C}_{\mathrm{AVE}}^{2}-1\right) / 2$ which has an asymptotic standard normal distribution (as $\mathrm{T} \rightarrow \infty$ first, then as $\mathrm{n} \rightarrow \infty$ ). However, when $\mathrm{n} \rightarrow \infty$ and T remains fixed, it turns out that the finite, and even the limiting, distribution of $\mathrm{C}_{\text {STD }}$ is not distribution-free. For example, under the null hypothesis of i.i.d. errors in (1.1), note that $\mathrm{EC}_{\mathrm{AVE}}^{2}=\mathrm{Ec}_{12}^{2}$ depends on the distribution of errors. See Section 2 for further discussion of this point. This dependence on the parent distribution suggests using a nonparametric version of $\mathrm{C}_{\mathrm{AVE}}^{2}$,

$$
\begin{equation*}
R_{A V E}^{2}=\left(\frac{n}{2}\right)^{-1} \Sigma_{i<j} r_{i j}^{2} \tag{1.4}
\end{equation*}
$$

where $r_{i j}$ is the Spearman rank correlation coefficient between the $i^{\text {th }}$ and $j^{\text {th }}$ flows. More specifically, define $r_{i j}$ $=\mathrm{s}_{\mathrm{ij}} /\left(\mathrm{s}_{\mathrm{ij}} \mathrm{s}_{\mathrm{j} j}\right)^{1 / 2}$, where $\mathrm{s}_{\mathrm{ij}}=(\mathrm{T}-1)^{-1} \Sigma_{\mathrm{t}}\left(\mathrm{R}_{\mathrm{i}, \mathrm{t}}-(\mathrm{T}+1) / 2\right)\left(\mathrm{R}_{\mathrm{j}, \mathrm{t}}-(\mathrm{T}+1) / 2\right)$ and $\left\{\mathrm{R}_{\mathrm{i}, \mathrm{l}}, \ldots, \mathrm{R}_{\mathrm{i}, \mathrm{T}}\right\}$ are the ranks of $\left\{\mathrm{Y}_{\mathrm{i} 1}\right.$, $\left.\ldots, \mathrm{Y}_{\mathrm{iT}}\right\}$. This statistic is distribution-free and the limiting distribution is established in Section 3. In applications where either positive or negative contemporaneous correlations prevail, one could also consider the statistic

$$
\begin{equation*}
R_{A V E}=\left(\frac{n}{2}\right)^{-1} \sum_{i<j} r_{i j} \tag{1.5}
\end{equation*}
$$

The advantage of this statistic is that the limiting distribution is known, that is, $(T-1)\left((n-1) R_{A V E}+1\right)=F R \rightarrow$
$\chi^{2}(T-1)$ where FR is Friedman's statistic, cf., Hetumansperger (1984, pp. 196, 210). For completeness, also define

$$
\begin{equation*}
C_{\mathrm{AVE}}=\binom{\mathrm{n}}{2}^{-1} \sum_{\mathrm{i}<\mathrm{j}} \mathrm{c}_{\mathrm{ij}} . \tag{1.6}
\end{equation*}
$$

to be the version of $\mathrm{R}_{\text {AVE }}$ using Pearson correlation coefficients.

In Section 2, I discuss the finite sample properties of the statistics in (1.3)-(1.6). In the case of $\mathrm{C}_{\mathrm{AVE}}^{2}$, this supplements the known asymptotic (for large T ) properties. The technique is to relate these statistics to a class of nonparametric unbiased estimators called U-statistics and use the well-known properties of this class. Because of the desirable properties of $\mathrm{R}_{\text {AVE }}^{2}$, in Section 3 the large (in n ) sample distribution of this statistic are established using classical techniques. Since the statistics are primarily for model diagnostics, they are most likely to be useful when evaluated using residuals from a preliminary model fit. Thus, in Section 4, I extend the discussion to statistics based on residuals. Sufficient conditions are provided so that the limiting distribution remains unchanged when using residuals in lieu of i.i.d. random variables. Further, the form of the distribution is identified when these conditions are not satisfied.

## 2. FINITE SAMPLE PROPERTIES

In this section, assume both n and T are fixed. Finite sample properties of the statistics (1.3)-(1.6) are developed using U-statistic theory, cf., Serfling (1980) for an introduction to these results.

Define $\mathrm{Y}_{\mathrm{i}}=\left(\mathrm{Y}_{\mathrm{il}}, \ldots, \mathrm{Y}_{\mathrm{If}}\right)^{\prime}$ and $\mathrm{e}_{\mathrm{i}}=\left(\mathrm{e}_{\mathrm{i} 1}, \ldots, \mathrm{e}_{\mathrm{It}}\right)^{\prime}$ to be random vectors in $\mathrm{R}^{\mathrm{T}}$, T-dimensional Euclidean space. Define four real-valued functions mapping $\mathrm{R}^{2 \mathrm{~T}}$ into R , called kernels, corresponding to the four statistics, as follows. Define $h_{1}\left(Y_{i}, Y_{j}\right)=c_{i j}, h_{2}\left(Y_{i}, Y_{j}\right)=r_{i j}, h_{3}=h_{1}^{2}$ and $h_{4}=h_{2}^{2}$. For each kernel, the $U$-statistic is defined to be the average over all possible evaluations of the kernel, that is,

$$
\begin{equation*}
U_{n, k}=\binom{n}{2}^{-1} \Sigma_{i<j} h_{k}\left(Y_{i}, Y_{j}\right) \quad k=1, \ldots, 4 . \tag{2.1}
\end{equation*}
$$

Note that $U_{n, 1}=C_{A V E}, U_{n, 2}=R_{A V E}, U_{n, 3}=C_{A V E}^{2}$ and $U_{n, 4}=R_{A V E}^{2}$. The simplicity of the results of this article rests in the fact that the correlation statistics are location and scale invariant so that $h_{k}\left(Y_{i}, Y_{j}\right)=h_{k}\left(e_{i}, e_{j}\right)$, that is, $h_{k}$ is a function of the i.i.d. random vectors $\left\{e_{i}\right\}$.

Throughout this paper, I assume

AI. The random variables $\left\{\mathrm{e}_{\mathrm{it}}\right\}$ are i.i.d.

This can be weakened at certain points to exchangeability and/or conditionally uncorrelated
( $E\left(h\left(e_{1}, e_{2}\right) \mid e_{1}\right)=0$ ), but for ease of interpretation I use the stronger assumption A1. It is straightforward to establish that $E U_{n, k}=E h_{k}$ for $k=1, \ldots, 4$. One can also check that $E R_{A V E}=E C_{A V E}=0$ under A1. The calculation of the variance is more complex. I summarize that calculation in the following

Theorem 1. Assume A1. Then

$$
\operatorname{Var} U_{n, k}=\binom{n}{2}^{-1} \quad \operatorname{Var} h_{k}\left(e_{1}, e_{2}\right) \quad k=1, \ldots, 4
$$

Before proving this result, it is helpful to note some important special cases.

COROLLARY 1. Assume A1 and that the random variables $\left\{\mathrm{e}_{\mathrm{it}}\right\}$ have a continuous distribution. Then $E R_{A V E}^{2}=(T-1)^{-1}, \operatorname{Var} R_{A V E}=\binom{n}{2}^{-1}(T-1)^{-1}$ and

$$
\operatorname{Var} R_{A V E}^{2}=\left(\frac{\mathrm{n}}{2}\right)^{-1}\left\{(\mathrm{~T}-2)\left(25 \mathrm{~T}^{2}-7 \mathrm{~T}-54\right) /\left(18(\mathrm{~T}-1)^{2}\left(\mathrm{~T}^{3}-\mathrm{T}\right)\right)\right\}
$$

Corollary 2. Assume Al and that the random variables $\left\{\mathrm{e}_{\mathrm{it}}\right\}$ are normally distributed. Then $E C_{A V E}^{2}=T^{-1}, \operatorname{Var} C_{A V E}=\left(\frac{n}{2}\right)^{-1} T^{-1}$ and

$$
\operatorname{Var} \mathrm{C}_{\mathrm{AVE}}^{2}=\binom{\mathrm{n}}{2}^{-1} 2(\mathrm{~T}-1) /\left\{\mathrm{T}^{2}(\mathrm{~T}+2)\right\}
$$

These corollaries provide explicit moments for fixed values of n and T . The assumption of a continuous distribution is the usual one in rank statistics, made to prevent possible ties in the random variables. The statistic $h_{1}^{2}\left(Y_{1}, Y_{2}\right)=h_{3}\left(Y_{1}, Y_{2}\right)$ is the squared correlation or coefficient of determination between $Y_{1}$ and $Y_{2}$. This is well-known to be depend on the distribution of $\left\{\mathrm{Y}_{\mathrm{it}}\right\}$, at least for finite T . In Corollary 2 is an example of the calculation in an important special case.

To prove Theorem 1, first recall some terminology from U-statistics theory. In this article, I consider only kernels that are symmetric in their arguments and are square integrable, that is, $h\left(y_{1}, y_{2}\right)=h\left(y_{2}, y_{1}\right)$ and $E h^{2}\left(Y_{1}, Y_{2}\right)<\infty$. Such a kernel is said to be degenerate if $\left.E\left(Y_{1}, Y_{2}\right) \mid Y_{1}=y_{1}\right)=E h\left(Y_{1}, Y_{2}\right)$ for all $y_{1}$. The
main work in establishing Theorem 1 is in

Lemma 1. Assume A1. Then $h_{k}$ is degenerate, for each $k=1, \ldots, 4$.

Proof of Lemma 1: It is straightforward to check that each $h_{k}$ is symmetric and square integrable. I only establish that $h_{4}$ is degenerate as the cases $h_{1}, h_{2}$ and $h_{3}$ are similar. For $h_{4}$, replace $Y$ by e in the definition of $\mathrm{b}_{\mathrm{ij}}$ so that $\mathrm{b}_{\mathrm{ij}}=(\mathrm{T}-1)^{-1} \Sigma_{\mathrm{t}}\left(\mathrm{e}_{\mathrm{it}}-\overline{\mathrm{e}}_{\mathrm{i}}\right)\left(\mathrm{e}_{\mathrm{jt}}-\overline{\mathrm{c}}_{\mathrm{j}}\right)$. Thus

$$
\begin{aligned}
E\left(h_{4}\left(e_{1}, e_{2}\right) \mid e_{2}\right)= & E\left(\left((T-1)^{-1} \Sigma_{\uparrow} e_{1 t}\left(e_{2 t}-\bar{e}_{2}\right)\right)^{2} /\left(b_{11} b_{22}\right) \mid e_{2}\right) \\
= & (T-1)^{-2} \Sigma_{t} \Sigma_{u}\left(e_{2 t}-\bar{e}_{2}\right)\left(e_{2 u}-\bar{e}_{2}\right) / b_{22} E\left(e_{1 t} e_{1 u} / b_{11}\right) \\
= & (T-1)^{-2} E\left(e_{11}^{2} / b_{11}\right) \Sigma_{\uparrow}\left(e_{2 t}-\bar{e}_{2}\right)^{2} / b_{22} \\
& \quad+E\left(e_{11} e_{12} / b_{11}\right) \Sigma_{\imath \neq u}\left(e_{2 t}-\bar{e}_{2}\right)\left(e_{2 u}-\bar{e}_{2}\right) / b_{22} \\
= & (T-1)^{-2}\left\{E\left(e_{11}^{2} / b_{11}\right)-E\left(e_{11} e_{12} / b_{11}\right)\right) \Sigma_{t}\left(e_{2 t}-\bar{e}_{2}\right)^{2} / b_{22} \\
= & (T-1)^{-1}\left\{E\left(e_{11}^{2} / b_{11}\right)-E\left(e_{11} e_{12} / b_{11}\right)\right\}
\end{aligned}
$$

which is a constant. Since $E E\left(h_{4}\left(e_{1}, e_{2}\right) \mid e_{1}\right)=E h_{4}\left(e_{1}, e_{2}\right)$, this constant must equal $E h_{4}$. This is sufficient for the result. Q. E. D.

The proof of Theorem 1 is an immediate consequence of Lemma 1 and the usual variance decomposition result for U-statistics (cf., Serfling, 1980, Lemma A, page 183). To establish Corollary 2, from standard linear model results, under normality we have $\mathrm{c}_{12}^{2}=\mathrm{F} /(\mathrm{F}+\mathrm{T}-1)$, where F has an F -distribution with 1 and T-1 degrees of freedom, respectively. Calculating the expected values of $c_{12}^{2}$ and $c_{12}^{4}$ is a pleasant exercise using moments of the $\chi^{2}$ - distribution and is omitted. Calculating moments of $r_{12}^{2}$ is also straightforward but, as some of these results are also used in Section 3, the details are sketched below.

Proof of Corollary 1: First note that, under the assumption of no ties, we have $\mathrm{s}_{\mathrm{ij}}$ is degenerate and equal to ( $\mathrm{T}^{2}$ + T)/12 (cf., Hetmansperger, 1984, p. 20). Now, without loss of generality, the ranks of one vector may be arranged in ascending order to get $h_{4}\left(e_{1}, e_{2}\right)=D \quad 144 /\left(T^{3}-T\right)^{2}\left(\sum_{t}(t-m)\left(R_{t}-m\right)\right)^{2}$ where $m=(T+1) / 2$ and $R_{t}$ is the rank of the $t^{\text {th }}$ element of the second vector. Basic calculations yield $E\left(R_{1}-m\right)^{2}=\left(T^{2}-1\right) / 12$ and $E\left(R_{1}-m\right)\left(R_{2}-m\right)=-(T+1) / 12$. This is sufficient to check that $E h_{4}\left(e_{1}, e_{2}\right)=1 /(T-1)$.

To calculate $E h_{4}^{2}$, the following quantities are useful: $S_{1}=\left(T^{3}-T\right) / 12$ and $S_{2}=\left(3 T^{2}-7\right) / 20$. Some basic equalities are $E\left(R_{1}-m\right)^{4}=S_{1} S_{2}, E\left(R_{1}-m\right)^{3}\left(R_{2}-m\right)=-S_{1} S_{2} /(T(T-1)), E\left(R_{1}-m\right)^{2}\left(R_{2}-m\right)\left(R_{3}-m\right)=S_{1}\left(2 S_{2}\right.$ $\left.-\mathrm{S}_{1}\right) /(\mathrm{T}(\mathrm{T}-1)(\mathrm{T}-2)), \mathrm{E}\left(\mathrm{R}_{1}-\mathrm{m}\right)^{2}\left(\mathrm{R}_{2}-\mathrm{m}\right)^{2}=\mathrm{S}_{1}\left(\mathrm{~S}_{1}-\mathrm{S}_{2}\right) /(\mathrm{T}(\mathrm{T}-1))$ and $\mathrm{E}\left(\mathrm{R}_{1}-\mathrm{m}\right)\left(\mathrm{R}_{2}-\mathrm{m}\right)\left(\mathrm{R}_{3}-\mathrm{m}\right)\left(\mathrm{R}_{4}-\mathrm{m}\right)=$ $3 \mathrm{~S}_{1}\left(\mathrm{~S}_{1}-2 \mathrm{~S}_{2}\right) /(\mathrm{T}(\mathrm{T}-1)(\mathrm{T}-2)(\mathrm{T}-3))$. After checking these basic equalities, the next major step is to verify that
$E h_{4}^{2}\left(c_{1}, c_{2}\right)=T\left(E\left(R_{1}-m\right)^{4}\right)^{2}+4 T(T-1)\left(E\left(R_{1}-m\right)^{3}\left(R_{2}-m\right)\right)^{2}+6 T(T-1)(T-2)\left(E\left(R_{1}-m\right)^{2}\left(R_{2}-m\right)\left(R_{3}-m\right)\right)^{2}+$ $3 \mathrm{~T}(\mathrm{~T}-1)\left(\mathrm{E}\left(\mathrm{R}_{1}-\mathrm{m}\right)^{2}\left(\mathrm{R}_{2}-\mathrm{m}\right)^{2}\right)^{2}+\mathrm{T}(\mathrm{T}-1)(\mathrm{T}-2)(\mathrm{T}-3)\left(\mathrm{E}\left(\mathrm{R}_{1}-\mathrm{m}\right)\left(\mathrm{R}_{2}-\mathrm{m}\right)\left(\mathrm{R}_{3}-\mathrm{m}\right)\left(\mathrm{R}_{4}-\mathrm{m}\right)\right)^{2}$.

Putting the basic equalities into this equation, and some tedious algebra, are sufficient for the result. Q. E. D.

## 3. LARGE SAMPLE PROPERTIES

I first provide a useful computational version of $R_{A V E}^{2}$ similar to one available for $R_{A V E}, c f$., Hettmansperger (1984, p.210). Define $Z_{i, t, u}=12\left(R_{i t}-m\right)\left(R_{i u}-m\right) /\left(T^{3}-T\right)$ where $m$ is the mean of $R_{i v}$, that is $\mathrm{m}=\mathrm{E} \mathrm{R}_{\mathrm{it}}=(\mathrm{T}+1) / 2$. Now, from (1.4), with $\mathrm{s}_{\mathrm{ii}}=\left(\mathrm{T}^{2}+\mathrm{T}\right) / 12$ and $\Sigma_{\mathrm{i}} \mathrm{Z}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}=0$, we have

$$
\begin{align*}
R_{A V E}^{2}= & 144 /\left(n(n-1)\left(T^{3}-T\right)^{2}\right) \Sigma_{i \neq j}\left(\sum_{t}\left(R_{i t}-m\right)\left(R_{j t}-m\right)\right)^{2} \\
& =(n(n-1))^{-1} \Sigma_{i \neq j} \Sigma_{t, u} z_{i, t, u} z_{j, t, u} \\
& =(n(n-1))^{-1} \Sigma_{t, u}\left(\left(\Sigma_{i} z_{i, t, u}\right)^{2}-\Sigma_{i} z_{i, t, u}^{2}\right) \tag{3.1}
\end{align*}
$$

For small T and large n , the expression of $\mathrm{R}_{\text {AVE }}^{2}$ in (3.1) is faster to compute than the definition in (1.4). The expression in (3.1) also provides some insights in the form of the limit distribution in Theorem 2. That is, the sums in (3.1) suggest limiting normal distributions, and when squared, suggest limiting weighted sums of $\chi^{2}$ distributions as appear in Theorem 2.

I now establish the following

Theorem 2. Assume A1 and that the random variables $\left\{\mathrm{e}_{\mathrm{it}}\right.$ \} have a continuous distribution. Then,

$$
\mathrm{n}\left(\mathrm{R}_{\mathrm{AVE}}^{2}-(\mathrm{T}-1)^{-1}\right) \rightarrow_{\mathrm{D}} \mathrm{Y}=\mathrm{a}(\mathrm{~T})\left(\chi_{1, \mathrm{~T}-1}^{2}-(\mathrm{T}-1)\right)+\mathrm{b}(\mathrm{~T})\left(\chi_{2, \mathrm{~T}(\mathrm{~T}-3) / 2}^{2}-\mathrm{T}(\mathrm{~T}-3) / 2\right)
$$

Here, $\chi_{1}^{2}$ and $\chi_{2}^{2}$ are independent $\chi^{2}$ random variables with $\mathrm{T}-1$ and $\mathrm{T}(\mathrm{T}-3) / 2$ degrees of freedom, respectively, and $\mathrm{a}(\mathrm{T})=4(\mathrm{~T}+2) /\left(5(\mathrm{~T}-1)^{2}(\mathrm{~T}+1)\right)$ and $\mathrm{b}(\mathrm{T})=2(5 \mathrm{~T}+6) /(5 \mathrm{~T}(\mathrm{~T}-1)(\mathrm{T}+1))$.

Remarks: Note that $\mathrm{E}=0$. Further, $\mathrm{a}(\mathrm{T})$ and $\mathrm{b}(\mathrm{T})$ are constants which depend on T such that $\mathrm{Var} \mathrm{Z}=$ $4(\mathrm{~T}-2)\left(25 \mathrm{~T}^{2}-7 \mathrm{~T}-54\right) /\left(25 \mathrm{~T}(\mathrm{~T}-1)^{3}(\mathrm{~T}+1)\right)=2 \operatorname{Var}_{4}\left(\mathrm{Y}_{1}, \mathrm{Y}_{2}\right)$. Since the limiting standardized asymptotic variance equals the variance of the limiting distribution, under uniform integrability, the fact that Var $\mathrm{Z} \sim \mathrm{n}$ Var $R_{\text {AVE }}^{2}$ is expected.

The method of proof below utilizes a conversion of the class of statistics to a quadratic form and then applies classical techniques. An alternative method of proof would have been to appeal to results on degenerate U-statistics with limiting distributions which are weighted sums of $\chi^{2}$, due to Gregory (1977) and Serfling (1980). However, these results were only explicitly stated for the case in which the arguments of the kernels are real-valued instead of vector-valued as is the case here. Presumably, this extension is minor but it seems that the classical techniques provides more insights into what makes the theorem work and computations needed to check the conditions underlying each technique are the same.

Proof of Theorem 2: Define the centered version of $Z, \dot{Z}_{i, t, u}=Z_{i, t, u}-E Z_{i, t, u}$ and note that $E Z_{i, t, u}=T^{-1}$ for $t=u$ and $=-1 /(T(T-1))$ for $t \neq u$. Starting with the second line of (3.1), with some algebra we have

$$
\begin{align*}
& n\left(R_{A V E}^{2}-(T-1)^{-1}\right)=n\left((n(n-1))^{-1} \sum_{i \neq j}\left\{\Sigma_{i, u} Z_{i, t, u} Z_{j, t, u}-(T-1)^{-1}\right\}\right. \\
& =(\mathrm{n}-1)^{-1} \Sigma_{\mathrm{i} \neq \mathrm{j}} \Sigma_{\mathrm{t}, \mathrm{u}} \dot{\mathrm{z}}_{\mathrm{i}, \mathrm{t}, \mathrm{u}} \dot{\mathrm{z}}_{\mathrm{j}, \mathrm{t}, \mathrm{u}} \\
& =(\mathrm{n}-1)^{-1} \Sigma_{\mathrm{t}, \mathrm{u}}\left(\left(\Sigma_{\mathrm{i}} \dot{\mathrm{Z}}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}\right)^{2}-\Sigma_{\mathrm{i}} \dot{\mathrm{Z}}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}^{2}\right) . \tag{3.2}
\end{align*}
$$

Now, by the weak law of large numbers, $(\mathrm{n}-1)^{-1} \Sigma_{\mathrm{i}} \dot{\mathrm{Z}}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}^{2} \rightarrow_{\mathrm{p}} \operatorname{Var}\left(\mathrm{Z}_{\mathrm{t}, \mathrm{u}}\right)$. Thus,

$$
\begin{gather*}
(n-1)^{-1} \quad \Sigma_{1,0} \quad \Sigma_{i} \dot{Z}_{i, L u}^{2}=T \operatorname{Var}\left(Z_{1,1}\right)+T(T-1) \operatorname{Var}\left(Z_{1,2}\right)+o_{p}(1) \\
=(T-2) /(T-1)+o_{p}(1) \tag{3.3}
\end{gather*}
$$

after some algebra similar to that in Corollary 2. Now define $Q_{n}$ to be a $T \times T$ matrix with $n^{-1 / 2} \Sigma_{i} \dot{Z}_{i, h, u}$ to be the element in the $t^{\text {th }}$ row and $u^{\text {th }}$ column. Further define $X_{n}=\operatorname{vec}\left(Q_{n}\right)$, the $T^{2} \times 1$ vector built by stacking the columns of $Q_{n}$. Then, using (3.2) and (3.3), we have

$$
\begin{equation*}
n\left(R_{A V E}^{2}-(T-1)^{-1}\right)=(1+o(1)) X_{n n}^{\prime}-(T-2) /(T-1)+o_{p}(1) \tag{3.4}
\end{equation*}
$$

By the multivariate central limit theorem, $\mathrm{X}_{\mathrm{n}} \rightarrow_{\mathrm{D}} \mathrm{X}={ }_{\mathrm{D}} \mathrm{N}(0, \Sigma)$, where $\Sigma=\mathrm{E}\left(\mathrm{X}_{1} \mathrm{X}_{1}^{\prime}\right)$. Putting this into (3.4), we have

$$
\mathrm{n}\left(\mathrm{R}_{\mathrm{AVE}}^{2}-(\mathrm{T}-1)^{-1}\right) \rightarrow_{\mathrm{D}} \mathrm{X}^{\prime} \mathrm{X}-(\mathrm{T}-2) /(\mathrm{T}-1)
$$

Thus, we only need to show that $Y={ }_{D} X^{\prime} X-(T-2) /(T-1)$. To this end, first compute the eigenvalues of $\Sigma$, defined by

$$
\begin{equation*}
\lambda \mathrm{a}=\Sigma \mathrm{a} . \tag{3.5}
\end{equation*}
$$

Here, ' $a$ ' is a $T^{2} \times 1$ vector built from $a=\operatorname{vec}(A)$ where $A$ is a $T \times T$ matrix with $a_{u}$ is the $t^{\text {th }}$ row and $u^{\text {th }}$ column. The $(t, u)^{\text {th }}$ component of (3.5) is

$$
\begin{equation*}
\lambda \mathrm{a}_{\mathrm{tu}}=\Sigma_{\mathrm{r}, \mathrm{~s}} \mathrm{a}_{\mathrm{rs}} \operatorname{Cov}\left(\mathrm{Z}_{\mathrm{tu}}, \mathrm{Z}_{\mathrm{rs}}\right)=\mathrm{E}\left(\dot{\mathrm{Z}}_{\mathrm{tu}} \Sigma_{\mathrm{rs}} \mathrm{a}_{\mathrm{rs}} \dot{\mathrm{Z}}_{\mathrm{rs}}\right) \tag{3.6}
\end{equation*}
$$

Recall $\Sigma_{\imath} \dot{z}_{t u}=0$ and note that $\Sigma_{\eta} \dot{z}_{\mathrm{t}}=0$. Thus, from (3.6), $\mathrm{a}_{\mathrm{tu}} \equiv 1$ yields $\lambda=0$ as an eigenvalue. To find non-zero eigenvalues, from (3.6), assuming $\lambda \neq 0$, we have $\Sigma_{u} a_{u t}=\Sigma_{u} a_{u u}=\Sigma_{u} a_{u u}=0$ for each $t$. With these constraints, now evaluate (3.6) for the case $t=u$. This yields

$$
\begin{aligned}
& \lambda \mathrm{a}_{\mathrm{tt}}= E\left(\dot{Z}_{\mathrm{t}} \Sigma_{\mathrm{rs}} \mathrm{a}_{\mathrm{rs}} \dot{\mathrm{Z}}_{\mathrm{rs}}\right) \\
&=\mathrm{a}_{\mathrm{tt}} \operatorname{Cov}\left(\dot{\mathrm{Z}}_{11}, \dot{\mathrm{Z}}_{11}\right)+\operatorname{Cov}\left(\dot{\mathrm{Z}}_{11}, \dot{\mathrm{Z}}_{12}\right)\left(\Sigma_{\mathrm{r}, \mathrm{~d}(\mathrm{r}, \mathrm{t})} \mathrm{a}_{\mathrm{r}}+\dot{a}_{\mathrm{tr}}\right) \\
&+\operatorname{Cov}\left(\dot{Z}_{11}, \dot{Z}_{22}\right)\left(\Sigma_{\mathrm{r}, \mathrm{~d}(\mathrm{r}, \mathrm{l})} a_{\mathrm{r}}\right)+\operatorname{Cov}\left(\dot{Z}_{11}, \mathrm{Z}_{23}\right)\left(\Sigma_{\mathrm{s}, \mathrm{r}, \mathrm{~d}(\mathrm{r}, \mathrm{~s}, \mathrm{t})} \mathrm{a}_{\mathrm{rs}}\right)
\end{aligned}
$$

$$
\begin{align*}
& =a_{t t}\left\{\operatorname{Cov}\left(\dot{Z}_{11}, \dot{Z}_{11}\right)-2 \operatorname{Cov}\left(\dot{Z}_{11}, \dot{Z}_{12}\right)-\operatorname{Cov}\left(\dot{Z}_{11}, \dot{Z}_{22}\right)-2 \operatorname{Cov}\left(\dot{Z}_{11}, \dot{Z}_{23}\right)\right\} \\
& =a_{t t} a(T) \tag{3.7}
\end{align*}
$$

Here, $\Sigma_{r, d(r, t)}$ means the sum over $\mathrm{r}=1, \ldots, \mathrm{~T}$ with distinct r and t . Similarly, $\Sigma_{\mathrm{s}, \mathrm{r}, \mathrm{d}(\mathrm{r}, \mathrm{s}, \mathrm{t})}$ means the sum over r and s with $\{\mathrm{r}, \mathrm{s}, \mathrm{t}\}$ distinct. Calculations similar to those in Corollary 2 show that the last equality in (3.7) holds. Now, for $t \neq u$,

$$
\begin{align*}
& \lambda \mathrm{a}_{\mathrm{tu}}=\mathrm{E}\left(\dot{\mathrm{Z}}_{\mathrm{tu}} \Sigma_{\mathrm{rs}} \mathrm{a}_{\mathrm{rs}} \dot{\mathrm{Z}}_{\mathrm{rs}}\right) \\
& =\operatorname{Cov}\left(\dot{Z}_{12}, \dot{Z}_{12}\right)\left(a_{t u}+a_{u t}\right)+\operatorname{Cov}\left(\dot{Z}_{11}, \dot{Z}_{12}\right)\left(a_{t t}+a_{u u}\right) \\
& +\operatorname{Cov}\left(\dot{Z}_{12}, \dot{Z}_{33}\right)\left(\Sigma_{r, d(r, u, t)} a_{r r}\right)+\operatorname{Cov}\left(\dot{Z}_{12}, \dot{Z}_{23}\right)\left(\sum_{r, d(r, u, t)} a_{r u}+a_{u r}+a_{t u}+a_{u t}\right) \\
& +\operatorname{Cov}\left(\dot{\mathrm{Z}}_{12}, \dot{\mathrm{Z}}_{34}\right) \Sigma_{\mathrm{r}, \mathrm{~s}, \mathrm{~d}(\mathrm{r}, \mathrm{~s}, \mathrm{t}, \mathrm{u})} \mathrm{a}_{\mathrm{rs}} \\
& =\left(\mathrm{a}_{\mathrm{tu}}+\mathrm{a}_{\mathrm{ut}}\right)\left\{\operatorname{Cov}\left(\dot{Z}_{12}, \dot{\mathrm{Z}}_{12}\right)-\operatorname{Cov}\left(\dot{\mathrm{Z}}_{12}, \dot{\mathrm{Z}}_{23}\right)+\operatorname{Cov}\left(\dot{\mathrm{Z}}_{12}, \dot{\mathrm{Z}}_{34}\right)\right\} \\
& +\left(\mathrm{a}_{11}+\mathrm{a}_{\mathrm{ut}}\right)\left\{\operatorname{Cov}\left(\dot{\mathrm{Z}}_{12}, \dot{\mathrm{Z}}_{11}\right)-\operatorname{Cov}\left(\dot{\mathrm{Z}}_{12}, \dot{\mathrm{Z}}_{33}\right)-\operatorname{Cov}\left(\dot{\mathrm{Z}}_{12}, \dot{\mathrm{Z}}_{23}\right)+2 \operatorname{Cov}\left(\dot{\mathrm{Z}}_{12}, \dot{\mathrm{Z}}_{34}\right)\right\} \\
& =\left(a_{u t}+a_{u t}\right) / 2 b(T)+\left(a_{t t}+a_{u u}\right) c(T) . \tag{3.8}
\end{align*}
$$

Here, $\Sigma_{r, s, d(r, s, t, u)}$ means the sum over $r$ and $s$ with $\{r, s, t, u\}$ distinct. Again, calculations similar to those in Corollary 2 yield the expressions $b(T)$ and $c(T)$. By (3.8), we have $a_{t u}=a_{u t}$.

For $\lambda \neq 0$, solutions to (3.7) - (3.8) are of two forms. First, if $a_{u t} \neq 0$ for some $t$, then $\lambda=a(T)$. Second, if $a_{t t}=0$ for each $t$, then $\lambda=b(T)$. Thus the eigenvalues of $\Sigma$ are either $0, a(T)$ or $b(T)$. It is also useful to recall that the sum of eigenvalues $=\operatorname{trace}(\Sigma)=\Sigma_{\mathrm{t}, \mathrm{u}} \operatorname{Var}\left(\mathrm{Z}_{\mathrm{tu}}\right)=(\mathrm{T}-2) /(\mathrm{T}-1)$, from (3.3). Now, there are T-1 free diagonal terms $a_{t t}$ since we have the restriction $\Sigma_{\tau} a_{t t}=0$. If any of these terms are non-zero, by (3.8), ( $a(T)$ $b(T)) a_{t u}=\left(a_{t t}+a_{u u}\right) c(T)$. Thus, fixing the diagonal terms of $A$ determines the remaining elements of the matrix and thus there are T-1 multiplicities at the root $\lambda=\mathrm{a}(\mathrm{T})$. If the diagonal terms of A are all zero, then there are T(T-1) off-diagonal terms to contend with. Due to the symmetry of A noted above and the restrictions $\Sigma_{u} a_{t u}=0$ for each $t$, it is casy to verify that there are $T(T-1) / 2-T=T(T-3) / 2$ free off-diagonal terms of $A$. This then establishes that there are $T(T-3) / 2$ multiplicities at root $\lambda=b(T)$. To establish that the only other roots are $\lambda=0$, we have that $(\mathrm{T}-1) \mathrm{a}(\mathrm{T})+\mathrm{T}(\mathrm{T}-3) / 2 \mathrm{~b}(\mathrm{~T})=\Sigma_{\mathrm{i}} \lambda_{\mathrm{i}}=(\mathrm{T}-2) /(\mathrm{T}-1)=\operatorname{trace}(\Sigma)$. That is, the sum of eigenvalues equals the trace, as it should.

The remainder of the proof is similar to the proof for $\mathrm{R}_{\mathrm{AVE}}$, cf., Hettmansperger (1984, page 184). There exists a $T^{2} \times T^{2}$ orthogonal matrix $G$ such that $G \Sigma G^{\prime}=D$, a diagonal matrix with the eigenvalues of $\Sigma$ on the diagonal. With $U=G X={ }_{D} N(0, D)$, we have

$$
\begin{aligned}
& \mathrm{X}^{\prime} \mathrm{X}-(\mathrm{T}-2) /(\mathrm{T}-1)=\mathrm{X}^{\prime} \mathrm{G}^{\prime} \mathrm{GX}-(\mathrm{T}-2) /(\mathrm{T}-1)=\mathrm{U}^{\prime} \mathrm{U}-(\mathrm{T}-2) /(\mathrm{T}-1) \\
& \quad=\sum_{\mathrm{i}} \lambda_{\mathrm{i}} \chi_{\mathrm{i}, 1}^{2}-(\mathrm{T}-2) /(\mathrm{T}-1)=\mathrm{a}(\mathrm{~T}) \chi_{1, \mathrm{~T}-1}^{2}+\mathrm{b}(\mathrm{~T}) \chi_{2, \mathrm{~T}(\mathrm{~T}-3) / 2}^{2}-(\mathrm{T}-2) /(\mathrm{T}-1) \\
& \quad={ }_{\mathrm{D}} \mathrm{Y} .
\end{aligned}
$$

Here, $\chi_{i, 1}^{2}$ are independent $\chi^{2}$ random variables with one degree of freedom. Q. E. D.

## 4. RESIDUAL BASED CORRELATIONS

As an extension of (1.1), now consider the model

$$
\begin{equation*}
Y_{i t}=\mu_{i}+g_{i t}(\theta)+\sigma_{i} e_{i t} \quad i=1, \ldots, n \quad t=1, \ldots, T \tag{4.1}
\end{equation*}
$$

where $g_{i t}$ is known up to $\theta$, a p x 1 vector of parameters. This is a nonlinear, heteroscedastic version of the model independently introduced by Hausman and Taylor (1981) and Laird and Ware (1982). Again, $\mu_{i}$ is a nuisance parameter and may be taken to be fixed or random. Let $\hat{\theta}$ be a root-n consistent estimate of $\theta$, i.e., assume

$$
\begin{equation*}
\mathrm{n}^{1 / 2}(\hat{\theta}-\theta)=\mathrm{O}_{\mathrm{p}}(1) \tag{4.2}
\end{equation*}
$$

Note that I do not require consistency of $\hat{\mu}_{i}$. Using these estimates and estimates of $\mu_{i}, \hat{\mu}_{i}$, define the residuals $\hat{e}_{i t}=Y_{i t}-\hat{Y}_{i t}=\mu_{i}-\hat{\mu}_{i}+g_{i t}(\theta)-g_{i t}(\hat{\theta})+\sigma_{i} e_{i t}$ and the corresponding ranks $\hat{R}_{i t}=\sum_{r=1}^{T} I\left(\hat{e}_{i r} \leq \hat{e}_{i t}\right)$. Analogous to (3.1), the contemporaneous correlation statistic based on residuals considered in this section is

$$
\begin{equation*}
\hat{R}_{A V E}^{2}=144 /\left(n(n-1)\left(T^{3}-T\right)^{2}\right) \sum_{i \neq j}\left(\sum_{\mathrm{i}}\left(\hat{R}_{i t}-m\right)\left(\hat{R}_{j t}-m\right)\right)^{2} \tag{4.3}
\end{equation*}
$$

Let $F($.$) and f($.$) be the distribution function and probability density function, assumed to exist, of the$ errors $\left\{\mathrm{e}_{\mathrm{it}}\right\}$. Some additional regularity conditions on the regression function $\mathrm{g}_{\mathrm{it}}$ are given in assumption A2
below. I am now in a position to state the main result of this section.

Theorem 3. Assume A1, A2, (4.2), $\inf _{\mathrm{i}} \sigma_{\mathrm{i}}>0$, and that the probability density function of $\left\{\mathrm{e}_{\mathrm{it}}\right\}$ exists. Then, under the model (4.1),

$$
n\left(R_{A V E}^{2}-\hat{R}_{A V E}^{2}\right)=\sum_{t=1}^{T}\left(W_{1, n}(t)+W_{2, n}(t)\right) W_{2, n}(t)+o_{p}(1)
$$

wherc $W_{1, n}(t)=2 n^{-1 / 2} \Sigma_{i} \dot{\mathrm{z}}_{\mathrm{i}, \mathrm{t}, \mathrm{l}}, \mathrm{W}_{2, \mathrm{n}}(\mathrm{t})=\mathrm{n}^{1 / 2}(\hat{\theta}-\theta)^{\prime}\left(\mathrm{n}^{-1} \Sigma_{\mathrm{i}} \nabla \mathrm{M}_{\mathrm{i}, \mathrm{l}}(\theta)\right)\left(12 /\left(\mathrm{T}^{3}-\mathrm{T}\right)\right)$ and $\nabla \mathrm{M}_{\mathrm{i}, \mathrm{t}}(\theta)$ is the gradient vector of $E_{\theta}\left(R_{t}(\theta)-m\right)^{2}$ which turns out to be

$$
\begin{equation*}
\nabla \mathrm{M}_{\mathrm{i}, \mathrm{t}}(\theta)=2(\mathrm{~T}-2)\left(\int(\mathrm{F}(\mathrm{x})-1 / 2) \mathrm{f}(\mathrm{x}) \mathrm{dF}(\mathrm{x})\right) / \sigma_{\mathrm{i}} \sum_{\mathrm{r}=1}^{\mathrm{T}}(\partial / \partial \gamma)\left[\mathrm{g}_{\mathrm{ir}}(\gamma)-\mathrm{g}_{\mathrm{it}}(\gamma)\right]_{\gamma=\theta} . \tag{4.4}
\end{equation*}
$$

Perhaps the most useful aspect of Theorem 2 is to give precise conditions when residuals behave as do i.i.d. error terms in the limiting distribution of $R_{A V E}^{2}$. An important special case is when $f($.$) is symmetric.$ Here, $\nabla \mathrm{M}_{\mathrm{i}, \mathrm{t}}(\theta)=0$ and we have

Corollary 1. Assume A1, A2, (4.2) and that random variables $\left\{\mathrm{e}_{\mathrm{it}}\right\}$ have a symmetric probability density function $\mathrm{f}($.$) . Then,$

$$
\mathrm{n}\left(\hat{\mathrm{R}}_{\mathrm{AVE}}^{2}-(\mathrm{T}-1)^{-1}\right) \rightarrow_{\mathrm{D}} \mathrm{Z}=\mathrm{a}(\mathrm{~T})\left(\chi_{1, \mathrm{~T}-1}^{2}-(\mathrm{T}-1)\right)+\mathrm{b}(\mathrm{~T})\left(\chi_{2, \mathrm{~T}(\mathrm{~T}-3) / 2}^{2}-\mathrm{T}(\mathrm{~T}-3) / 2\right)
$$

The method of analysis employs the "estimated parameter" approach advocated by Randles (1984). To this end, I use Randles' (1984, Condition 1.10) regularity condition on the regression function, as follows.

A2. Let $K(\theta)$ be a neighborhood of $\theta$ in $R^{p}$. For each $\gamma \in K(\theta)$, assume that $\mathrm{g}_{\mathrm{it}}(\gamma)$ is differentiable uniformly in i and the gradient vector, $\partial / \partial \gamma \mathrm{g}_{\mathrm{it}}(\gamma)$ is uniformly bounded in i .

Implicit in A 2 , as noted by Randles, is that if $\mathrm{g}_{\mathrm{it}}$ contains covariates then the covariate space is bounded. Now, for $\gamma \in R^{p}$, define the perturbed errors, $\mathrm{e}_{\mathrm{it}}(\gamma)=\mathrm{g}_{\mathrm{it}}(\theta)-\mathrm{g}_{\mathrm{it}}(\gamma)+\sigma_{\mathrm{i}} \mathrm{e}_{\mathrm{it}}$, so that $\mathrm{e}_{\mathrm{it}}(\hat{\theta})=\hat{\mathrm{e}}_{\mathrm{it}}$ up to $\mu_{\mathrm{i}}-\hat{\mu}_{\mathrm{i}}$ which, being merely location shifts, do not affect the subsequent analysis. We also need the perturbed ranks

$$
\begin{equation*}
R_{i t}(\gamma)=\sum_{r=1}^{T} I\left(e_{i r}(\gamma) \leq e_{i t}(\gamma)\right)=\sum_{r=1}^{T} I\left(e_{i r} \leq e_{i t}+\eta_{i, r, t}(\gamma)\right) \tag{4.5}
\end{equation*}
$$

where $\eta_{i, r, t}(\gamma)=\left\{g_{i t}(\theta)-g_{i t}(\gamma)-g_{i r}(\theta)+g_{i r}(\gamma)\right\} / \sigma_{i}$. Note that $\eta_{i, r, t}(\gamma)=-\eta_{i, t, r}(\gamma)$. Similar to Section 3, define $Z_{i, \mathrm{~L}, \mathrm{u}}(\gamma)=12\left(\mathrm{R}_{\mathrm{it}}(\gamma)-\mathrm{m}\right)\left(\mathrm{R}_{\mathrm{iu}}(\gamma)-\mathrm{m}\right) /\left(\mathrm{T}^{3}-\mathrm{T}\right)$. The main work in proving Theorem 3 is in

## LEMMA 2. Under the assumptions of Theorem 3, for each $t, u$,

$$
\mathrm{n}^{-1 / 2}\left(\Sigma_{\mathrm{i}} \mathrm{Z}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}(\hat{\theta})-\mathrm{Z}_{\mathrm{i}, \mathrm{i}, \mathrm{u}}(\theta)\right)-\mathrm{I}(\mathrm{t}=\mathrm{u}) \mathrm{W}_{2, \mathrm{n}}(\mathrm{t}) \rightarrow_{\mathrm{p}} 0
$$

Proof of Lemma 2: The lemma is a consequence of Theorem A. 9 of Randles' (1984). To verify the conditions of that theorem, first note that (4.2) and Assumption A2 satisfy Randles' A. 4 and Condition 1.10, respectively. Define the kernel of order 1 to be

$$
h_{i}\left(e_{i}, \gamma\right)=Z_{i, r, u}(\gamma)
$$

It is easy to see that the kernel is uniformly bounded (satisfying Randles' A.8a). Now, following Randles' notation, for $K(\theta)$ in $A 2$, let $D(\gamma, d) \in R^{p}$ be a sphere centered at $\gamma$ with radius such that $D$ is contained in $K(\theta)$. We wish to prove that there exists $K>0$ so that

$$
\begin{equation*}
\mathrm{q}_{\mathrm{i}}(\mathrm{~d}) \equiv \mathrm{E}_{\theta} \sup _{\gamma^{\prime} \in \mathrm{D}}\left|\mathrm{~h}_{\mathrm{i}}\left(\mathrm{e}_{\mathrm{i}}, \gamma\right)-\mathrm{h}_{\mathrm{i}}\left(\mathrm{e}_{\mathrm{i}}, \gamma^{\prime}\right)\right| \leq \mathrm{Kd} \tag{4.6}
\end{equation*}
$$

By A2, there exists $K_{1}>0$ such that $\sup _{\gamma^{\prime} \in D}\left|g_{i t}(\gamma)-g_{i t}\left(\gamma^{\prime}\right)\right| \leq K_{1}$ d. This, and since inf $\sigma_{i}>0$, implies that there exists $K_{2}>0$ such that $\sup _{\gamma^{\prime} \in D}\left|\eta_{i, r, t}(\gamma)-\eta_{i, r, t}\left(\gamma^{\prime}\right)\right| \leq K_{2}$ d. From this and (4.5),

$$
\begin{aligned}
& E_{\theta} \sup _{\gamma^{\prime} \in D}\left|R_{i t}(\gamma)-R_{i t}\left(\gamma^{\prime}\right) I=E_{\theta} \sup _{\gamma^{\prime} \in D}\right| \sum_{r} I\left(e_{i r} \leq e_{i t}+\eta_{i, r, t}(\gamma)\right)-I\left(e_{i r} \leq e_{i t}+\eta_{i, r, t}\left(\gamma^{\prime}\right)\right) \mid \\
& \quad \leq E_{\theta} \sup _{\gamma^{\prime} \in D} \sum_{r, d(r, t)} I\left(\left|e_{i r}-e_{i t}\right| \leq\left|\eta_{i, r, t}(\gamma)-\eta_{i, r, t}\left(\gamma^{\prime}\right)\right|\right) \\
& \quad \leq \sum_{r, d(r, t)} E_{\theta} I\left(e_{i r}-e_{i t} \mid \leq K_{2} d\right) \\
& \quad \leq(T-2) K_{2} d .
\end{aligned}
$$

Recall that $\sum_{r, d(r, t)}$ means the sum over $r$ with distinct $\{r, t\}$. This and the triangle inequality are sufficient for (4.6).

We only need now calculate the gradient of $E_{\theta} h_{i}\left(e_{i}, \gamma\right)=M_{i, u}(\gamma)$ and check that it is uniformly (in i) achieved at $\gamma=\theta$. First consider the case $\mathrm{t}=\mathrm{u}$. From (4.5), we have

$$
\begin{equation*}
\mathrm{R}_{\mathrm{it}}(\gamma)-\mathrm{m}=\sum_{\mathrm{r}, \mathrm{~d}(\mathrm{r}, \mathrm{t})}\left(\mathrm{I}\left(\mathrm{e}_{\mathrm{ir}} \leq \mathrm{e}_{\mathrm{it}}+\eta_{\mathrm{i}, \mathrm{r}, \mathrm{t}}(\gamma)\right)-1 / 2\right) \tag{4.7}
\end{equation*}
$$

Thus,

$$
\begin{align*}
M_{i, t}(\gamma) \equiv & E_{\theta}\left(R_{i t}(\gamma)-m\right)^{2} \\
= & E_{\theta} \Sigma_{r, d(r, t)}\left(I\left(c_{i r} \leq c_{i t}+\eta_{i, r, t}(\gamma)\right)-1 / 2\right)^{2} \\
& \quad+E_{\theta} \sum_{r, s, d,(r, s, t)}\left(I\left(e_{i r} \leq e_{i t}+\eta_{i, r, t}(\gamma)\right)-1 / 2\right)\left(I\left(e_{i s} \leq e_{i t}+\eta_{i, s, t}(\gamma)\right)-1 / 2\right) \\
= & (T-1) / 4+\Sigma_{r, s, d,(r, s, t)} E_{\theta}\left(F\left(e+\eta_{i, r, t}(\gamma)\right)-1 / 2\right)\left(F\left(e+\eta_{i, s, t}(\gamma)\right)-1 / 2\right) \tag{4.8}
\end{align*}
$$

Note that $(\partial / \partial \gamma) \eta_{\mathrm{i}, \mathrm{r}, \mathrm{t}}(\gamma)=(\partial / \partial \gamma)\left[\mathrm{g}_{\mathrm{ir}}(\gamma)-\mathrm{g}_{\mathrm{it}}(\gamma)\right] / \sigma_{\mathrm{i}}$. By the chain rule, for fixed $\mathrm{e}, \mathrm{i}, \mathrm{r}, \mathrm{s}$, and u , we have

$$
\begin{align*}
& (\partial / \partial \gamma)\left[\left(\mathrm{F}\left(\mathrm{e}+\eta_{\mathrm{i}, \mathrm{r}, \mathrm{t}}(\gamma)\right)-1 / 2\right)\left(\mathrm{F}\left(\mathrm{e}+\eta_{\mathrm{i}, \mathrm{~s}, \mathrm{l}}(\gamma)\right)-1 / 2\right)\right]_{\gamma=0} \\
& \quad=(\mathrm{F}(\mathrm{e})-1 / 2) \mathrm{f}(\mathrm{e}) / \sigma_{\mathrm{i}}(\partial / \partial \gamma)\left[\mathrm{g}_{\mathrm{ir}}(\gamma)+\mathrm{g}_{\mathrm{is}}(\gamma)-2 \mathrm{~g}_{\mathrm{it}}(\gamma)\right]_{\gamma=\theta} \tag{4.9}
\end{align*}
$$

Now take the vector derivative of (4.8). Since $r$ and $s$ are finite, an application of the Bounded Convergence Theorem, with A2 and (4.9), yields $\nabla \mathrm{M}_{\mathrm{i}, \mathrm{t}}(\boldsymbol{\theta})$ in (4.4). Showing that this is achieved uniformly in i is another application of the Bounded Convergence Theorem.

For $t \neq u$, first recall that $\eta_{i, u, t}=-\eta_{i, t, u}$ and thus $I\left(e_{i u} \leq e_{i t}+\eta_{i, u, t}(\gamma)\right)-1 / 2=1 / 2-I\left(e_{i t} \leq e_{i u}+\eta_{i, t, u}(\gamma)\right)$.
Thus, from (4.7)

$$
\begin{aligned}
& E_{\theta}\left(R_{i t}(\gamma)-m\right)\left(R_{i u}(\gamma)-m\right) \\
& =\mathrm{E}_{\theta}\left[\sum_{\mathrm{r}, \mathrm{~d}(\mathrm{r}, \mathrm{t}, \mathrm{u})}\left(\mathrm{I}\left(\mathrm{e}_{\mathrm{ir}} \leq \mathrm{e}_{\mathrm{it}}+\eta_{\mathrm{i}, \mathrm{r}, \mathrm{t}}(\gamma)\right)-\mathrm{l} / 2\right)+\left(\mathrm{I}\left(\mathrm{e}_{\mathrm{iu}} \leq \mathrm{e}_{\mathrm{it}}+\eta_{\mathrm{i}, \mathrm{u}, \mathrm{t}}(\gamma)\right)-1 / 2\right)\right] \\
& {\left[\sum_{r, d,(r, t, u)}\left(I\left(e_{i r} \leq e_{i u}+\eta_{i, r, u}(\gamma)\right)-1 / 2\right)+\left(I\left(e_{i t} \leq e_{i u}+\eta_{i, t, u}(\gamma)\right)-1 / 2\right)\right]} \\
& =E_{\theta}\left\{\left[\sum_{r, d(r, t, u)}\left(I\left(e_{i r} \leq e_{i t}+\eta_{i, r, t}(\gamma)\right)-1 / 2\right)\right]\left[\Sigma_{r, d,(r, t, u)}\left(I\left(e_{i r} \leq e_{i u}+\eta_{i, r, u}(\gamma)\right)-1 / 2\right)\right]\right. \\
& +\left[\left(1\left(e_{i u} \leq \mathrm{e}_{\mathrm{it}}+\eta_{\mathrm{i}, \mathrm{u}, \mathrm{t}}(\gamma)\right)-1 / 2\right)\right] \\
& {\left[\Sigma_{r, d,(r, t, u)}\left(I\left(e_{i r} \leq e_{i t}+\eta_{i, r, t}(\gamma)\right)-I\left(e_{i r} \leq e_{i u}+\eta_{i, r, u}(\gamma)\right)\right]\right\}-1 / 4} \\
& \left.=\sum_{r, d(r, t, u)} E_{\theta}\left(I\left(e_{i r} \leq e_{i t}+\eta_{i, r, i}(\gamma)\right)-1 / 2\right)\left(I\left(e_{i r} \leq e_{i u}+\eta_{i, r, u}(\gamma)\right)-1 / 2\right)\right] \\
& +\sum_{r, s, d,(r, s, t)} E_{\theta}\left(F\left(e+\eta_{i, r, i}(\gamma)\right)-1 / 2\right) E_{\theta}\left(F\left(e+\eta_{i, r, u}(\gamma)\right)-1 / 2\right) \\
& +\sum_{r, d(r, t, u)} E_{\theta}\left(I\left(e_{i u} \leq e_{i t}+\eta_{i, u, t}(\gamma)\right)-1 / 2\right)\left[F\left(e_{i t}+\eta_{i, r, t}(\gamma)\right)-\left(F\left(e_{i u}+\eta_{i, r, u}(\gamma)\right)\right]-1 / 4 .\right.
\end{aligned}
$$

Since $E_{\theta}(F(e)-1 / 2)=0$, by the chain rule, we have

$$
\begin{aligned}
& (\partial / \partial \gamma)\left[\mathrm{E}_{\theta}\left(\mathrm{R}_{\mathrm{it}}(\gamma)-\mathrm{m}\right)\left(\mathrm{R}_{\mathrm{iu}}(\gamma)-\mathrm{m}\right)\right]_{\gamma=\theta} \\
& =(\partial / \partial \gamma)\left[\sum_{r, d(r, r, u)} E_{\theta}\left(\mathrm{I}\left(\mathrm{e}_{\mathrm{ir}} \leq \mathrm{e}_{\mathrm{it}}+\eta_{\mathrm{i}, \mathrm{r}, \mathrm{i}}(\gamma)\right)-1 / 2\right)\left(\mathrm{I}\left(\mathrm{e}_{\mathrm{ir}} \leq \mathrm{e}_{\mathrm{iu}}+\eta_{\mathrm{i}, \mathrm{r}, \mathrm{u}}(\gamma)\right)-1 / 2\right)\right. \\
& +\sum_{r, d(r, t, u)} E_{\theta}\left(I\left(e_{i u} \leq e_{i t}+\eta_{i, u, t}(\gamma)\right)-1 / 2\right)\left(F\left(e_{i t}+\eta_{i, r, t}(\gamma)\right)-\left(F\left(e_{i u}+\eta_{i, r, u}(\gamma)\right)\right)\right]_{\gamma=\theta} \\
& =(\partial / \partial \gamma)\left[\sum _ { \mathrm { r } , \mathrm { d } ( \mathrm { r } , \mathrm { t } , \mathrm { u } ) } \left[\mathrm { E } _ { \theta } \left\{\left(\mathrm{F}\left(\mathrm{e}+\eta_{\mathrm{i}, \mathrm{t}, \mathrm{r}}(\gamma)\right)-1 / 2\right)\left(\mathrm{F}\left(\mathrm{e}+\eta_{\mathrm{i}, \mathrm{u}, \mathrm{r}}(\gamma)\right)-1 / 2\right)\right.\right.\right. \\
& +\left(\mathrm{F}\left(\mathrm{e}+\eta_{\mathrm{i}, \mathrm{u}, \mathrm{t}}(\gamma)\right)-1 / 2\right)\left(\mathrm{F}\left(\mathrm{e}+\eta_{\mathrm{i}, \mathrm{r}, \mathrm{t}}(\gamma)\right)-1 / 2\right) \\
& \left.\left.+\left(F\left(e+\eta_{i, t, u}(\gamma)\right)-1 / 2\right)\left(F\left(e+\eta_{i, r, u}(\gamma)\right)-1 / 2\right)\right\}\right]_{\gamma=\theta} \\
& =\left(\mathrm{E}_{\theta}(\mathrm{F}(\mathrm{e})-1 / 2) \mathrm{f}(\mathrm{e}) / \sigma_{\mathrm{i}}\right)(\partial / \partial \gamma)\left[\mathrm{g}_{\mathrm{il}}(\gamma)+\mathrm{g}_{\mathrm{iu}}(\gamma)-2 \mathrm{~g}_{\mathrm{ir}}(\gamma)\right. \\
& \left.+\mathrm{g}_{\mathrm{iu}}(\gamma)+\mathrm{g}_{\mathrm{ir}}(\gamma)-2 \mathrm{~g}_{\mathrm{it}}(\gamma)+\mathrm{g}_{\mathrm{it}}(\gamma)+\mathrm{g}_{\mathrm{ir}}(\gamma)-2 \mathrm{~g}_{\mathrm{iu}}(\gamma)\right]_{\gamma=\theta} \\
& =0,
\end{aligned}
$$

since $(\partial / \partial \gamma) E_{\theta}\left[\left(\mathrm{F}\left(\mathrm{e}+\eta_{\mathrm{i}, \mathrm{u}, \mathrm{i}}(\gamma)\right)+\mathrm{F}\left(\mathrm{e}+\eta_{\mathrm{i}, \mathrm{t}, \mathrm{u}}(\gamma)\right)\right]_{\gamma=\theta}=0\right.$. This establishes that the gradient of $\mathrm{E} \mathrm{Z}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}(\gamma)$ at $\gamma=\theta$ is zero. The fact that it is uniformly achieved in ' $i$ ' can be established similarly to the case $t=u$. Q. E. D.

## Proof of Theorem 3: From (4.3) and (3.1), we have

$$
n\left(R_{A V E}^{2}-\hat{R}_{A V E}^{2}\right)=(1+o(1)) n^{-1} \Sigma_{t, u}\left\{\left(\Sigma_{i} \dot{Z}_{i, t, u}\right)^{2}-\left(\sum_{i} \dot{Z}_{i, t, u}(\hat{\theta})\right)^{2}-\left(\Sigma_{i} \dot{Z}_{i, t, u}{ }^{2}-\dot{Z}_{i, t, u}(\hat{\theta})^{2}\right)\right\}
$$

where $\dot{Z}_{i, t, u}(\hat{\theta}) \equiv \mathrm{Z}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}(\hat{\theta})-\mathrm{E} \mathrm{Z}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}(\theta)$. By Lemma $2, \mathrm{n}^{-1} \sum_{\mathrm{i}}\left(\dot{\mathrm{Z}}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}-\dot{\mathrm{Z}}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}(\hat{\theta})\right) \rightarrow_{\mathrm{p}} 0$. Similar arguments can be used to show that $n^{-1} \sum_{i}\left(\dot{Z}_{i, t, u}{ }^{2}-\dot{Z}_{i, t, u}(\hat{\theta})^{2}\right) \rightarrow_{p} 0$. Thus,

$$
\begin{aligned}
\mathrm{n}\left(\mathrm{R}_{\mathrm{AVE}}^{2}-\hat{\mathrm{R}}_{\mathrm{AVE}}^{2}\right) & =\mathrm{n}^{-1} \Sigma_{\mathrm{t}, \mathrm{u}}\left\{\left(\Sigma_{\mathrm{i}} \dot{Z}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}-\dot{Z}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}(\hat{\theta})\right)\left(\Sigma_{\mathrm{i}} \dot{Z}_{\mathrm{i}, \mathrm{t}, \mathrm{u}}+\dot{Z}_{\mathrm{i}, \mathrm{,u}}(\hat{\theta})\right)\right\}+\mathrm{o}_{\mathrm{p}}(1) \\
& =\Sigma_{\mathrm{u}}\left(\mathrm{I}(\mathrm{t}=\mathrm{u}) \mathrm{W}_{2 \mathrm{n}}(\mathrm{t})+\mathrm{o}_{\mathrm{p}}(1)\right)\left(\mathrm{W}_{1 \mathrm{n}}(\mathrm{t})+\mathrm{W}_{2 \mathrm{n}}(\mathrm{t})+\mathrm{o}_{\mathrm{p}}(1)\right)+\mathrm{o}_{\mathrm{p}}(1)
\end{aligned}
$$

which, with A2, (4.2) and Lemma 2, is sufficient for the proof. Q. E. D.

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