On the Meaning of Certain Cointegration Tests
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# ON THE MEANING OF CERTAIN COINTEGRATION TESTS* 

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#### Abstract

In this paper we examine (some of) the cointegration tests suggested in Engle and Granger (1987), henceforth EG, those suggested by Johansen (1988), (1991), henceforth J, and those suggested in Dhrymes (1994b), henceforth D. We also explore the relations among the models underlying the various procedures. We find that the tests suggested by Johansen cannot possibly test either for the presence of cointegration or for its rank, but they can determine the cointegration vectors, if cointegration is known to hold and its rank is known as well.


Key words: Integrated Processes, Stationary Processes, Collinearity, Cointegration, Cointegration tests.

## 1 Introduction and Preliminaries

This paper narrowly focuses on some of the suggestions in EG, J and the tests in D. In Dickey, Jansen and Thornton (1991) (DJT) there is a good review of the procedures in EG and $J$, so that we shall not focus on them in great detail unless, as in the case of $J$, the argument warrants it.

In all three approaches the basic problem is that we are confronted with a sequence $\left\{X_{t}: t \geq 1\right\}$ which is known to be $I(1)$, and where $X_{t}$. is a $q$-element row vector. The question posed is whether or not it is

[^0]cointegrated. For a definition of cointegration, rank thereof etc. see, for example, Dhrymes (1994). The underlying models are
\[

$$
\begin{align*}
& (I-L) X_{t .}^{\prime}=A(L) \epsilon_{t .}^{\prime}, \quad A(L)=\sum_{j=0}^{k} A_{j} L^{j}, \quad \text { Dhrymes } \\
& (I-L) X_{t}^{\prime} .=C(L) \epsilon_{t .}^{\prime}, \quad C(L)=\sum_{j=0}^{\infty} C_{j} L^{j}, \quad \text { Engle-Granger } \\
& \Pi(L) X_{t .}^{\prime}=\epsilon_{t .}^{\prime}, \quad \Pi(L)=\sum_{j=1}^{p} \Pi_{j} L^{j}, \quad \text { Johansen. } \tag{1}
\end{align*}
$$
\]

In the first two, the $\epsilon$-sequence is taken to be a $M W N(\Sigma)$, i.e. a multivariate white noise process with mean zero and covariance matrix $\Sigma>0$. In either case if the argument required it, the authors would not be averse to asserting that higher moments existed as well. In the third case, J , the $\epsilon$-sequence is asserted to be a Gaussian $M W N(\Sigma)$.

The first two formulations are essentially the same, in practical terms, since in $D$ the order of the moving average, $k$, need not be known a priori. Moreover, it is shown therein that the EG formulation is not very apt, given the problem to be investigated. If it is taken seriously, i.e. if it is applied verbatim in conjunction with the usual definition of cointegration, it cannot admit of cointegration. See the examination of the example given in EG, which makes this point in Dhrymes (1994a). Essentially, the problem is the presence of the requisite initial conditions, in the absence of which no such discussions may proceed. The nature of the test in D, however, may be viewed as in the spirit of the original EG intent, if not necessarily the original exposition. The third (J) formulation is somewhat different, in that it neither contains nor is contained, as a special case, in the other two formulations.

## 2 Models and Implied Tests

In the first formulation (D), a test for cointegration is immediately available once we establish that

$$
\begin{aligned}
\operatorname{Cov}\left(X_{t .}^{\prime}\right) & =\sum_{r=0}^{t-1} S_{0, r} \Sigma S_{0, r}^{\prime}, \quad \text { if } t \leq k \\
& =(t-k) S_{0, k} \Sigma S_{0, k}^{\prime}+\sum_{r=0}^{k-1} S_{0, r} \Sigma S_{0, r}^{\prime}, \text { if } t \geq k+1, \text { or }
\end{aligned}
$$

$$
\begin{align*}
& =t \Phi+D, \quad \Phi=S_{0, k} \Sigma S_{0, k}^{\prime} \\
D & =\sum_{r=0}^{k-1}\left(S_{0, k} \Sigma S_{0, k}^{\prime}-S_{0, r} \Sigma S_{0, r}^{\prime}\right), \tag{2}
\end{align*}
$$

where

$$
\begin{equation*}
S_{i, j}=\sum_{s=i}^{j} A_{s}, \quad i, j=0,1,2, \ldots, k, \quad S_{0,0}=A_{0}=I_{q} \tag{3}
\end{equation*}
$$

Since, by the definition of cointegration, a necessary and sufficient (nas) condition for cointegration of rank $r$, is the existence of a $q \times r$ matrix $B$ of rank $r$, such that the cointegral vector, $Z_{t}=X_{t} . B$, has a stationary covariance matrix (or more precisely is covariance stationary), it is clear that $B$ must be in the null space of $S_{0, k} \Sigma S_{0, k}^{\prime}$ but not in that of $D$, and the latter must of rank at least $r$. Hence, in this formulation the (a) set of cointegrating vectors is simply the basis of the null space of $S_{0, k} \Sigma S_{0, k}^{\prime}$. Since $\Sigma>0$, we may also state that the (a) set of cointegrating vectors is simply the basis of the null space of $S_{0, k}$. Thus, in this framework we have cointegration if and only if there is at least one zero among the roots of $S_{0, k} \Sigma S_{0, k}^{\prime}$ and the cointegrating rank is simply the number of such zero roots.

The instrumentality of carrying out this test is the estimator

$$
\begin{equation*}
M_{T,(1)}=\frac{1}{T} \sum_{t=1}^{T} \frac{X_{t}^{\prime} X_{t} .}{t}, \quad \text { or, alternatively, } \quad M_{T,(2)}=\frac{1}{T^{2}} \sum_{j=1}^{T} X_{t}^{\prime} X_{t} \tag{4}
\end{equation*}
$$

The first is an asymptotically unbiased estimator of the desired matrix, while the second is an asymptotically biased estimator of $\Phi$, the limit of its expectation being ( $1 / 2$ ) $\Phi$. The problem with both estimators is that under the hypothesis that the $X$-sequence is $I(1)$, their limiting distribution is nonstandard and needs to be established and tabulated.

In the D context to test whether the $X$-process is $I(0)$ is equivalent to testing whether all roots of $\Phi$ are null. An alternative recommended in Dhrymes (1994b) is to actually test for zero roots in

$$
\begin{equation*}
M_{T,(0)}=\frac{1}{T} \sum_{t=1}^{T} X_{t}^{\prime} X_{t} \tag{5}
\end{equation*}
$$

under the hypothesis that the $X$-process is $I(0)$. The argument here is that if zero roots are detected in this context, it would mean that we have collinearity among the underlying variables, and whatever else is obtained by a proper cointegration test may simply be a reflection of this collinearity, and the well known phenomenon noted first by Stone (1947), and discussed in the author's early work, Dhrymes (1970). The
desirability of this procedure, as a way of ruling out contaminated results is strengthened by

Proposition 1. Consider the matrices $M_{T,(i)}, i=0,1$, of the previous discussion; the matrix $M_{T,(0)}-M_{T,(1)}$, conditionally on the sample, is positive semidefinite.

Proof: Neglecting the factor $(1 / T)$, we find

$$
X^{\prime} X-X^{\prime} N^{2} X=X^{\prime} D X, \quad N=\operatorname{diag}\left(\frac{1}{t^{1 / 2}}\right), \quad D=\operatorname{diag}\left(\frac{t-1}{t}\right)
$$

for $t=1,2, \ldots, T$. Since $D \geq 0$, it follows that $X^{\prime} D X \geq 0$.
q.e.d.

Corollary 1. If, conditionally on the sample, $M_{T,(0)}$ has $r_{0}$ roots that obey $\mu_{j}^{(0)} \leq \delta$ then there exist $r_{1} \geq r_{0}$ characteristic roots of $M_{T,(1)}$, say $\mu_{s}^{(1)}$ such that $\mu_{s}^{(1)} \leq \delta, s=1,2 \ldots r_{1} \geq r_{0}$.

Proof: We note, Bellman (1960) p. 115, that since

$$
\begin{equation*}
X^{\prime} X=X^{\prime} N^{2} X+X^{\prime} D X \tag{6}
\end{equation*}
$$

and both matrices on the right are at least positive semidefinite, it follows that

$$
\begin{equation*}
\mu_{j}^{0} \geq \mu_{j}^{(1)}, \quad j=1,2,3, \ldots, q \tag{7}
\end{equation*}
$$

It follows, therefore, that

$$
\begin{equation*}
\mu_{j}^{(1)} \leq \mu_{j}^{(0)} \leq \delta, \quad j=1,2,3, \ldots, r_{0} \tag{8}
\end{equation*}
$$

q.e.d.

Remark 1. Evidently, the implications to be derived from Proposition 1 and Corollary 1 above are suggestive, not conclusive.

We now consider the model put forth by J. In this formulation we begin with a $M A R(p)$ model, generally known in macroeconomics as a VAR, i.e.

$$
\begin{equation*}
\Pi(L) X_{t}^{\prime}=\epsilon_{t .}^{\prime}, \quad \Pi(L)=\sum_{j=0}^{p} \Pi_{j} L^{j}, \quad \Pi_{0}=I_{q} \tag{9}
\end{equation*}
$$

where the $\epsilon$-process is a $M W N(\Sigma)$ Gaussian process. After considerable manipulation, but without any transformation of the error process or of the underlying model, we may write

$$
\begin{equation*}
(I-L) X_{t .}^{\prime}=-\Pi(I) X_{t-1 \cdot}^{\prime}-\sum_{j=1}^{p} \Pi_{j}^{*}(I-L) X_{t-j .}^{\prime}+\epsilon_{t \cdot}^{\prime} \tag{10}
\end{equation*}
$$

where, of course, $\Pi(I)$ is the operator $\left(\sum_{i=0}^{p} \Pi_{i}\right) I$. In current paractice, the literature does not distinguish between the operator $\Pi(I)$, and the matrix $\Pi(I)=\sum_{i=0}^{p} \Pi_{i}$. Following this custom, we shall not distinguish between the two in subsequenct discussion. ${ }^{1}$ It is important to understand that the conceptual basis of this approach is different from those in the D or EG formulation, and we shall return to this issue below. For the moment, let us note that if the series is cointegrated of rank $r$ then $\Pi(I)$ must be a $q \times q$ matrix of rank $r$. This is so since the left and right members of Eq. (10) are all $I(0)$, except possibly for the first term of the right member. Since we assert that the series $X$ is cointegrated of rank $r$ then the term $\Pi(I) X_{t}^{\prime}$. must be, at most, a multiple of the cointegral vector $Z_{t-1 .}=X_{t-1} \cdot B$, where $B$ is a $q \times r$ matrix of rank $r$, whose columns are the cointegrating vectors. Thus, by the singular value decomposition theorem, we may write

$$
\begin{equation*}
\Pi(I)=\Gamma B^{\prime}, \text { where } \Gamma, B \text { are (both) } q \times r \text { matrices of rank } r \tag{11}
\end{equation*}
$$

It bears stressing that in order to write $\Pi(I)$ as in Eq. (11) we must assert cointegration of rank $r$. Otherwise $\Pi(I)$ is just another parameter matrix, whose rank can be determined by a test performed utilizing the estimated parameter (or a function thereof) as a test statistic. Under the maintained hypothesis that we have cointegration of rank $r$, the J procedure maximizes the LF with respect to the ancillary parameters $\Sigma, \Pi_{j}^{*}, j=1,2, \ldots p$, and $\Gamma$, thus obtaining the concentrated LF

$$
\begin{equation*}
L^{*}(B)=-\frac{q T}{2}[\ln (2 \pi)+1]-\frac{T}{2} \ln |S|, \quad S=\frac{1}{T} W^{\prime} N^{*} W . \tag{12}
\end{equation*}
$$

To justify the representation in Eq. (12), define

$$
\begin{align*}
y_{t .} & =(I-L) X_{t .}-X_{t-1} \cdot \Pi(I)^{\prime}=\Delta X_{t .}-X_{t-1} \cdot B \Gamma^{\prime} \\
x_{t .} & =\left(\Delta X_{t-1} \cdot \Delta X_{t-2}, \Delta X_{t-3 \cdot}, \ldots, \Delta X_{t-n+1} \cdot\right), \quad \epsilon_{t}=u_{t \cdot} \\
\Pi^{*} & =\left(-\Pi_{1}^{*^{\prime}},-\Pi_{2}^{*^{\prime}}, \ldots, \Pi_{n-1}^{*^{\prime}}\right)^{\prime}, \quad W=N \Delta P, \quad V=N P_{-1} \\
Y & =\left(y_{t \cdot}\right), \quad X=\left(x_{t \cdot}\right), \quad U=\left(u_{t \cdot} \cdot\right), t=1,2,3, \ldots T \\
P & =\left(X_{t .}\right), \quad P_{-i}=\left(X_{t-i} \cdot\right), \quad t=1,2,3, \ldots, T, \quad i=1,2,3, \ldots n-1 \\
p & =\operatorname{vec}(P), \quad p_{-i}=\operatorname{vec}\left(P_{-i}\right), \quad N=I_{T}-X\left(X^{\prime} X\right)^{-1} X^{\prime} \\
N^{*} & =I_{T}-N P_{-1} B\left(B^{\prime} P_{-1}^{\prime} N P_{-1} B\right)^{-1} B^{\prime} P_{-1}^{\prime} N,  \tag{13}\\
y & =\operatorname{vec}(Y), \quad \pi^{*}=\operatorname{vec}\left(\Pi^{*}\right), \quad u=\operatorname{vec}(U), \quad \gamma=\operatorname{vec}\left(\Gamma^{\prime}\right)
\end{align*}
$$

[^1]so that we can write
\[

$$
\begin{equation*}
y_{t .}=x_{t} . \Pi^{*}+u_{t .}, \quad Y=X \Pi^{*}+U, \quad \text { or } \quad y=\left(I_{q} \otimes X\right) \pi^{*}+u \tag{14}
\end{equation*}
$$

\]

for a single observation and the entire sample, respectively. The concentrated LF is obtained by partially maximizing the loglikelihood (LF) of the $T$ observations

$$
\begin{equation*}
L_{0}=-\frac{q T}{2} \ln (2 \pi)-\frac{T}{2} \ln |\Sigma|-\frac{1}{2} \operatorname{tr}\left(Y-X \Pi^{*}\right) \Sigma^{-1}\left(Y-X \Pi^{*}\right)^{\prime} \tag{15}
\end{equation*}
$$

From Dhrymes (1984) p. 106, we have
$\operatorname{tr}\left(Y-X \Pi^{*}\right) \Sigma^{-1}\left(Y-X \Pi^{*}\right)^{\prime}=\left[y-\left(I_{q} \otimes X\right) \pi^{*}\right]^{\prime}\left(\Sigma^{-1} \otimes I_{T}\right)\left[y-\left(I_{q} \otimes X\right) \pi^{*}\right]$, and solving the first order condition $\left(\partial L_{0} / \partial \pi^{*}\right)=0$, we find

$$
\begin{equation*}
\hat{\pi^{*}}=\left[I_{q} \otimes\left(X^{\prime} X\right)^{-1} X^{\prime}\right] y \tag{16}
\end{equation*}
$$

Inserting this in Eq. (12), we obtain the concentrated LF

$$
\begin{equation*}
L_{\mathbf{1}}=-\frac{q T}{2} \ln (2 \pi)-\frac{T}{2} \ln |\Sigma|-\frac{1}{2} y^{\prime}\left(\Sigma^{-1} \otimes N\right) y \tag{17}
\end{equation*}
$$

Next, we note that

$$
\begin{equation*}
y=\Delta p-\left(I_{q} \otimes P_{-1} B\right) \gamma \tag{18}
\end{equation*}
$$

and the concentrated LF may be rendered as

$$
\begin{align*}
L_{1}=- & \frac{q T}{2} \ln (2 \pi)-\frac{T}{2} \ln |\Sigma|  \tag{19}\\
& -\frac{1}{2}\left[\Delta p-\left(I_{q} \otimes P_{-1} B\right) \gamma\right]^{\prime}\left(\Sigma^{-1} \otimes N\right)\left[\Delta p-\left(I_{q} \otimes P_{-1} B\right) \gamma\right]
\end{align*}
$$

solving the first order conditions $\left(\partial L_{1} / \partial \gamma\right)=0$, we find

$$
\begin{equation*}
\hat{\gamma}=\left[I_{q} \otimes\left(B^{\prime} P_{-1}^{\prime} N P_{-1} B\right)^{-1} B^{\prime} P_{-1}^{\prime} N\right] \Delta p \tag{20}
\end{equation*}
$$

Noting that

$$
\begin{equation*}
\left(I_{q} \otimes N\right)\left[\Delta p-\left(I_{q} \otimes P_{-1} B\right) \hat{\gamma}\right]=\left(I_{q} \otimes N^{*}\right)\left(I_{q} \otimes N\right) \Delta p \tag{21}
\end{equation*}
$$

we may write the (once again) concentrated LF as

$$
\begin{align*}
L_{2}=- & \frac{q T}{2} \ln (2 \pi)-\frac{T}{2} \ln |\Sigma|  \tag{22}\\
& -\frac{1}{2}\left[\left(I_{q} \otimes N^{*} N\right) \Delta p\right]^{\prime}\left(\Sigma^{-1} \otimes I_{T}\right)\left[\left(I_{q} \otimes N^{*} N\right) \Delta p\right]
\end{align*}
$$

Using the results in Dhrymes (1984) p. 106 and "rematricizing" the last expression in the LF, we may rewrite the latter as

$$
\begin{equation*}
L_{2}=-\frac{q T}{2} \ln (2 \pi)-\frac{T}{2} \ln |\Sigma|-\frac{T}{2} \operatorname{tr} \Sigma^{-1} S, \quad S=\frac{1}{T}(N \Delta P)^{\prime} N^{*}(N \Delta P) . \tag{23}
\end{equation*}
$$

Maximizing $L_{2}$ with respect to the elements of $\Sigma^{-1}$, we obtain

$$
\begin{equation*}
\frac{\partial L_{2}}{\partial \operatorname{vec}\left(\Sigma^{-1}\right)}=\frac{T}{2} \operatorname{vec}(\Sigma)^{\prime}-\frac{T}{2} \operatorname{vec}(S)=0, \quad \text { or } \quad \hat{\Sigma}=S \tag{24}
\end{equation*}
$$

Inserting this in Eq. (23), we find the ultimately concentrated LF, $L^{*}$, as it appears in Eq. (12). This is now to be maximized with respect to $B$, i.e. we are to find a $B$ that maximizes the LF. It bears stressing again that, in the absence of prior restrictions, to complete the estimation procedure implicit in the J formulation nothing more and nothing less is required by that formulation than simply to maximize the function in Eq. (12) with respect to the elements of the matrix $B$. For the moment, let us proceed with this unrestricted framework and see where it leads us. Without loss of generality, and for notational simplicity, we neglect the factor $(1 / T)$ in the definition of $S$ and, equivalently, pose the problem: minimize $S$ with respect to $B$. Now, put

$$
D(B)=T^{q}|S|=\left|W^{\prime} W-W^{\prime} V\left(W^{\prime} W\right)^{-1} V^{\prime} W\right|
$$

and note that, after some manipulation, we obtain

$$
\begin{align*}
D(B) & =\left|W^{\prime} W\right|\left|B^{\prime} V^{\prime} V B\right|^{-1}\left|B^{\prime}(F-H) B\right| \\
F & =V^{\prime} V, \quad H=V^{\prime} W\left(W^{\prime} W\right)^{-1} W^{\prime} V \tag{25}
\end{align*}
$$

Since $D(B)$ is homogeneous of degree zero in $B$, the problem is not well defined, and a normalization is required. A convenient normalization in this context is $B^{\prime} V^{\prime} V B=I_{q}$. It is a rather straightforward procedure to show that the minimizing value of $B$ is given by the matrix of the characteristic vectors of $H$ in the metric of $F .^{2}$ Thus, the

[^2]minimum value of $D(B)$, obtained without any prior restrictions, is given by
\[

$$
\begin{equation*}
\min _{B^{\prime} V^{\prime} V B=I_{q}} D(B)=\left|W^{\prime} W\right| \prod_{j=1}^{q}\left(1-\lambda_{j}\right), \tag{26}
\end{equation*}
$$

\]

where the $\lambda_{j}$ are the characteristic roots of $H$, in the metric of $F$. The question that arises now, in view of the fact that we are dealing with the loglikelihood function, relates to the magnitude of these roots. For if some roots are greater than one the maximum of the LF is not well defined, and thus the problem has no solution. To investigate this issue consider the matrix

$$
G=\left(P_{-1}, \Delta P\right)^{\prime} N\left(P_{-1}, \Delta P\right)=\left[\begin{array}{cc}
V^{\prime} V & V^{\prime} W \\
W^{\prime} V & W^{\prime} W
\end{array}\right]
$$

and note that since neither $P_{-1}$, nor $\Delta P$ is in the null space of $N$, $G \geq 0$. Defining

$$
K=\left[\begin{array}{cc}
I & -V^{\prime} W\left(W^{\prime} W\right)^{-1} \\
0 & I
\end{array}\right]
$$

we conclude from

$$
K G K^{\prime}=\left[\begin{array}{cc}
F-H & 0  \tag{27}\\
0 & W^{\prime} W
\end{array}\right]
$$

that $F-H=V^{\prime} V-V^{\prime} W\left(W^{\prime} W\right)^{-1} W^{\prime} V$ is also positive semidefinite. Hence, the characteristic roots of $H$ in the metric of $F$, obey $\lambda_{j} \in$ $[0,1)$. If the matrix $F-H$ is singular then at least one of the $\lambda_{j}$ is equal to one, and hence the maximum of the LF is not well defined. ${ }^{3}$ Zero roots, if any, are due to the singularity of $H$. Thus, provided no root is unity, the LF maximized without prior restrictions yields

$$
\begin{equation*}
\max _{B, \Sigma, \Gamma^{*}, \Gamma} L_{3}=-\frac{q T}{2} \ln [(2 \pi)+1]-\frac{1}{T} \ln \left|W^{\prime} W\right|-\frac{T}{2} \sum_{j=1}^{q} \ln \left(1-\lambda_{j}\right) . \tag{28}
\end{equation*}
$$

Since $1-\lambda_{j}>0$, if there are zero roots their inclusion or exclusion is a matter of complete indifference from the point of view of estimation since the LF is flat relative to this eventuality. Thus, whether to choose characteristic vectors corresponding to zero roots is akin to the situation encountered in many nonlinear problems where the likelihood function

[^3]may have more than one maxima. Strictly speaking this results in lack of identifiction! Thus, the problem as posed by J does not admit of a solution, unless restrictions are imposed a priori. In fact the only unambiguous result that may be obtained from this procedure, absent any prior restrictions, is in the case where it is conlcuded that $B$ is nonsingular, without the occasion of zero roots!

But even if we arrive at the nonsingularity of $B$ without the intervention of zero roots, we still face a problem of degeneracy. To see this, return to the definition of $S$ and note that if $B$ is nonsingular, the determinant $|S|$ does not involve $B$, hence $B$ is arbitrary subject to its being nonsingular. Moreover, returning to Eq. (20) and noting the nonsingularity of $B$ we obtain, after rematricizing the expression for $\hat{\Gamma}$,

$$
\begin{equation*}
\hat{\Pi}(I)=\hat{\Gamma}^{\prime} B=\left(P_{-1}^{\prime} N P_{-1}\right)^{-1} P_{-1}^{\prime} N \Delta P \tag{29}
\end{equation*}
$$

whose right member has the appearance of a "2SLS-like" estimator of $\Pi(I)$ in the pseudo-relation $\Delta P=N P_{-1} \Pi(I)+$ error.

If now we impose the prior restrictions (a) that the series is cointegrated and (b) cointegration is of rank $r<q$, we have a well defined problem and an unambiguous solution, $B_{r}$, even if one of the roots, whose corresponding characteristic vector is a column of $B_{\tau}$, is zero. The reason why this is so is because in this context, we have no idea what cointegration implies beyond the requirement that $\Pi(I)$ be singular and of rank $r<q$, and the procedure has, as required, produced an estimator of $\Pi(I)$, viz. $\hat{\Gamma} \hat{B}_{r}^{\prime}$, which is a $q \times q$ matrix of rank $r$, irrespective of whether some of the characteristic roots in question are null. Thus, the "admissibility" of null characteristic roots is not an issue in this problem. The difficulty is that, absent prior information on the rank and existence of cointegration, the estimation problem is not well defined, and thus cannot yield relevant tests about the existence of cointegration or the magnitude of its rank!

Remark 2. It was the practice, until the advent of the J procedure, for empirical researchers to carry out a Dickey-Fuller (DF) test, or an "augmented" DF test, and then proceed to test for cointegration and estimate the underlying parameters. The problem was that the several tests suggested in Engle and Granger (1987) (EG), and subsequent modifications that produced a number of variants, were either too cumbersome, or not very well argued, or their sampling properties were not available. For whatever reason, the current practice is to employ almost exclusively the J procedure. This is aided by the fact that the relevant test statistics involve the chi-squared distribution, which is well tabulated. An additional
factor, no doubt, is that the last stage of maximization procedure, i.e. the minimization of $|S|$ with respect to $B$ involves certain mechanics which are almost identical to some part of the process of obtaining pairs of canonical variates. No doubt this similarity adds to the intuitive appeal of the J procedure, even though this similarity is coincidental, and completely irrelevant. The canonical variate problem is the following: given two sets of (zero mean) random vectors, say $x^{1}$ and $x^{2}$, and each of dimension $q$, find $r$ pairs of linear combinations of $x^{1}$ and $x^{2}$, respectively, that exhibit maximal correlation, subject to their being pairwise uncorrelated (with other pairs of linear combinations), each element of the pair having unit variance. If the linear combinations in question are given by $G_{r}^{\prime} x^{1}$ and $H_{r}^{\prime} x^{2}$, respectively, their constituent vectors, $g_{\cdot i}, h_{\cdot i}, i=1,2, \ldots, r$, are the solutions to the system

$$
\left[\begin{array}{ll}
\lambda \Sigma_{11} & -\Sigma_{12}  \tag{30}\\
-\Sigma_{21} & \lambda \Sigma_{22}
\end{array}\right]\binom{g}{h}=0
$$

which leads to

$$
\begin{equation*}
\left(\lambda^{2} \Sigma_{11}-\Sigma_{12} \Sigma_{22}^{-1} \Sigma_{21}\right) g=0, \quad h=\frac{1}{\lambda} \Sigma_{22}^{-1} \Sigma_{21} g \tag{31}
\end{equation*}
$$

The first equation in Eq. (31) defines the vectors $g_{. i}$ as the characteristic vectors corresponding to the characteristic roots $\lambda_{i}^{2}$ of the appropriate determinantal equation. The second equation defines the vectors $h_{. i}$. Moreover, it can be shown that the correlation between the components of each pair may be displayed in the diagonal matrix, $\Lambda_{r}$, below, i.e.

$$
\begin{equation*}
E G_{r}^{\prime} x^{1} x^{2^{\prime}} H_{r}=\Lambda_{r}=\operatorname{diag}\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{r}\right) \tag{32}
\end{equation*}
$$

it being understood that $\lambda_{i}$ is the correlation between $g_{\cdot i}^{\prime} x^{1}$ and $h_{. i}^{\prime} x^{2}$. It is only the part associated with the first equation in Eq. (31) that resembles the mechanics of the J procedure. In the case of canonical variates one rejects vectors corresponding to zero roots, either because of the second equation in Eq. (31) above, which would involves division by zero, or because the roots of the first equation may be interpreted as the square of the correlations between the two components of a pair of canonical variates. In the $J$ procedure, however, there is no reason inherent in the formulation of the problem that may lead one to reject characteristic vectors corresponding to zero roots, particularly since the operation involved in the second equation of Eq. (31) does not enter the $J$ procedure. The latter is not designed to, nor does it obtain the canonical variates to which we may transform $N P_{-1}$ and $N \Delta P$. It would be difficult to visualize such an operation unless one were to argue that $N P_{-1}$ represents observations on a stationary
process. Thus, the fact that (part of) the estimation mechanics of the J procedure coincide in with (part of) the mechanics of obtaining canonical variates is completely coincidental and irrelevant to the operation, or interpretation of the results obtained in the J context. As noted above, zero roots present identification problems. As in the simultaneous equations literature identification problems may be "solved" only by means of prior restrictions. There is no more justification in rejecting all zero roots than in rejecting only some of them. For example, if in the J context there are $s>1$ roots we may admit $s_{0}<s$, thereby establishing cointegration of rank $q-\left(s-s_{0}\right)$, or admit all of them, thus establishing that the series is $I(0)$, or admit none of them in which case we shall conclude that we have cointegration of rank $q-s$. From the point of view of the maximum of the LF there is absolutely no distinction among these alternatives. If we claim that we know that cointegration is present, then we are entitled to take the rank of $B$ to be $q-1$, since this would make $\Pi(I)$ of rank $q-1$. Equally validly, we could take its rank to be one! Unless we know what the rank is we cannot settle on anything but maximal rank, which implies degeneracy. It is important to bear in mind that whether one or more of the relevant roots employed is zero has nothing to do with the question of the rank of the matrix $B$ we have chosen! Thus, testing to see whether some characteristic roots are, or are not equal to zero is very much like dueling with windmills-it does not affect the outcome of anything of any consequence.

The next question of significance is: why are the results from J so much different from those in EG and Dhrymes (1994b). To answer this question let us juxtapose the three models, obtaining

$$
\begin{align*}
D & :(I-L) X_{t .}^{\prime}=A(L) \epsilon_{t .}^{\prime}, \quad A(L)=\sum_{j=0}^{k} A_{j} L^{j}, \quad A_{0}=I_{q} \\
E G: & (I-L) X_{t .}^{\prime}=C(L) \epsilon_{t .}^{\prime}, \quad C(L)=\sum_{j=0}^{\infty} C_{j} L^{j}, \quad \sum_{j=0}^{\infty}\left\|C_{j}\right\|<\infty \\
& C_{0}=I_{q}  \tag{33}\\
J & : \Pi(L) X_{t .}^{\prime}=\epsilon_{t .}^{\prime}, \quad \Pi_{0}=I_{q} .
\end{align*}
$$

The first two models are basically similar, with the exception that the ostensible greater generality of the EG formulation leads to difficulties with the concept and existence of cointegration in such models. On the other hand, any application of models that derive from the EG formu-
lation will have to rely on a finite number of parameters only, which is the point of the D formulation to begin with. Thus, we may view D as a refinement and correction of the EG formulation in so far as the comparison with the $J$ formulation is concerned. In the first two formulations we may write, without loss of generality

$$
\begin{equation*}
A(L)=A(I)+(I-L) A^{*}(L), \quad C(L)=C(I)+(I-L) C^{*}(L) \tag{34}
\end{equation*}
$$

which are obtained by the long division of the operators $A(L)$ and $C(L)$, respectively by $-(L-I)$, coupled in the case of $C(L)$ with a limiting argument. In Eq. (34) $A(I), C(I)$ are the remainders and $A^{*}(L), C^{*}(L)$ the quotients. In the EG case further conditions must be satisfied by $C(L)$ is order for $C^{*}(L)$ to be well defined. A similar decomposition may be made of $\Pi(L)$, thus obtaining

$$
\begin{equation*}
\Pi(L)=\Pi(I)+(I-L) \Pi^{*}(L) . \tag{35}
\end{equation*}
$$

The first thing to observe in all three formulations is the role played by the remainder terms. If $A(I)=0$, or $C(I)=0$, the D and EG formulation suggest that the $X$-process is $I(0)$. If $\Pi(I)=0$, on the other hand, the J formulation suggests that the $X$-process is $I(1)$, provided $\left|\Pi^{*}(z)\right|=0$, does not have a root on the unit circle, i.e. if $z_{j}$ is any root then $\left|z_{j}\right| \neq 1$, more precisely $\left|z_{j}\right|<1 .{ }^{4}$ We shall see in a moment what is the consequence of having $\left|\Pi^{*}\left(z_{0}\right)\right|=0$, for $\left|z_{0}\right|=$ 1. If $A(I), C(I)$ are nonsingular, we cannot have cointegration and, in the case of the D formulation at least, $A(L)$ is an invertible operator. ${ }^{5}$ In either case, the nonsingularity of the remainder matrix, allows the $X$-process to remain $I(1)$ wihout being cointegrated. If $A(I), C(I) \neq 0$, but are (both) singular, we have cointegration and the rank of cointegration is simply the dimension of the null space of these matrices, respectively for the D and EG formulations. Similarly, in the case of the J formulation if $\Pi(I)$ is nonsingular, $\Pi(L)$ is invertible, for the same reasons given above, and we conclude that the $X$-process is $I(0)$. Finally if $\Pi(I) \neq 0$, but it is singular, we have cointegration and the rank of cointegration is the dimension of the null space of $\Pi(I)$.

Notice now certain important distinctions:

[^4]i. $A(I)=0$ implies the $X$-sequence is $I(0)$ in the D formulation; $\Pi(I)=0$ implies the $X$-sequence is $I(1)$ in the J formulation. Notice that, in an inference context, it is almost impossible to arrive at these conclusions with real world data in the $J$ formulation; in the D formulation this is possible if and only if the process is, indeed, $I(0)$ and the sample is sufficiently large;
ii. $|A(I)|=0$ implies cointegration in the D context, and similarly $|\Pi(I)|=0$ implies cointegration in the J formulation and, in either, the rank of cointegration is the dimension of the null space of the matrices corresponding to $A(I)$ and $\Pi(I)$, respectively;
iii. $A(I)$ and $\Pi(I)$ nonsingular (and implicitly assuming there are no complex roots of unit modulus) imply in the $D$ formulation that the $X$-process is $I(1)$ without cointegration, while in the J process it implies that the $X$-process is $I(0)$.

Notice that the two formulations are sort of polar opposites, and this is to be expected since one relies on a $M M A(k)$ representation of the differenced $X$-process, while the other relies on an $\operatorname{MAR}(p)$ representation of the $X$-process. The range of alternatives in the D formulation is: $I(0), I(1)$ with cointegration, $I(1)$ without cointegration. For the J formulation it is: $I(0)$ and $I(1)$ with cointegration. The alternative $I(1)$ without cointegration is almost impossible to attain since, in terms of the earlier discussion would entail something like $P_{-1}^{\prime} N \Delta P=0$, or more precisely a "positive outcome" to the test of the hypothesis that $B=0$. The reader should note that in the J estimation context, there is no way in which $B$ could be estimated to be zero, or even insignificantly different from zero.

Remark 3. If one wished to invent a valid test in the $J$ formulation, one of the possible approaches may be: estimate consistently $\Pi(I)$, perhaps by least squares, determine the distribution of $\hat{\Pi}(I)^{\prime} \hat{\Pi}(I)$ and test for the nummber of nonzero (positive) roots of this matrix; if the number so determined is $r$, this is the cointegration rank. If no roots are determined to be positive one would conclude that $\Pi(I)=0$ and the $X$-process is $I(1)$, without cointegration; if no root is determined to be zero one would conclude that $\Pi(I)$ is nonsingular, and thus the $X$-sequence is $I(0)$. This will restore to the test the full range of outcomes that a proper test for cointegration should have.

Remark 4. What is it that is so deficient in the J formulation that it leads to such undesirable features? The basic reason is that, except in
the strange case that one writes $\Pi(L)$ when one means to write ( $I-$ $L) \Pi^{*}(L)$, the formulation does not permit the $X$-process to be $I(1)$ without cointegration. Note that in the $D$ or $E G$ formulation the characteristic equation of the system of difference equations is

$$
\begin{equation*}
0=\left|(1-z) I_{q}\right|=(1-z)^{q} \tag{36}
\end{equation*}
$$

so that the unit root in question is of multiplicity $q$. In the J formulation the characteristic equation is

$$
\begin{equation*}
|\Pi(z)|=(1-z)\left|\Pi_{(1)}(z)\right|=0 \tag{37}
\end{equation*}
$$

if $\Pi(z)$ has only one unit root; if it has a unit root of multiplicity $r<q$, it may be represented as

$$
\begin{equation*}
|\Pi(z)|=(1-z)^{r}\left|\Pi_{(r)}(z)\right| \tag{38}
\end{equation*}
$$

and the reader should note that $\left|\Pi_{(1)}(z)\right|$ is of degree at most $q p-1$, while the determinant in Eq. (38) is a polynomial of degree at most $q p-r$. Thus, the multiplicity of degree $r$ version of the J formulation implies that the long run behavior of the $X$-process is essentially that of

$$
\Pi_{(r)}(L)(I-L)^{r / q} X_{t .}^{\prime}=K(L) \epsilon_{t \cdot}^{\prime}
$$

for an appropriate forcing function, $K(L) \epsilon_{t .}^{\prime}$, which is $M M A(s)$ with $s<\infty$. While the relationship above is only suggestive, and need not imply that cointegration is synonymous with fractional integration, nonetheless it illuminates the question of why and how the J formulation is inappropriate for testing for the presence of, or for the rank of, cointegration.

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    $\dagger$ This is a preliminary verision and is not to be quoted, except by permission of the author. Comments, however, are welcome.

[^1]:    ${ }^{1}$ Note that in our notation we distringuish between the identity operator $I$ and the $q \times q$ identity matrix $I_{q}$.

[^2]:    ${ }^{2}$ This may be illustrated quite simply when $B$ is taken to be a single vector, say $b$. The problem then is to maximize $b^{\prime} H b$ subject to $b^{\prime} F b=1$. Set up the Lagrangian $\mathcal{L}=b^{\prime} H b+\lambda\left(1-b^{\prime} F b\right)$. The first order conditions yield $(\lambda F-H) b=0$, and the solution is evidently given as one of the characteristic vectors of $H$ in the metric of $F$; since we wish to maximize, we choose as the solution the vector, say $b^{*}$, that corresponds to the largest root, $\lambda^{*}$. Thus we obtain, for the orginal problem,

    $$
    b^{*^{\prime}}(F-H) b^{*}=1-\lambda^{*},
    $$

    which yields, for $b^{\prime}(F-H) b$, the minimum we wish. Since the equation that defines this solution is the characteristic equation of $H$ in the metric of $F$, which has $q$ roots and corresponding characteristic vectors, the reader would easily grasp the conclusion that in a problem that involves all such roots and vectors, such as the minimization of a determinant, the solution would involve all characteristic roots and vectors as well.

[^3]:    ${ }^{3}$ The argument in this connection is to be understood as an a.c. (almost certain) argument. If one wished, one could convert all such entities to their parametric equivalents, if they are ascertained to exist. There is no difficulty with $W^{\prime} W$ which involves stationary entities only. There may be a problem with $V^{\prime} V$, since $V=$ $N P_{-1}$, and $P_{-1}$ involves elements that may not be stationary; on the other hand $V^{\prime} V / T^{2}$, as well as $H / T^{2}$ converge, at least in distribution, and the a.c. argument may be made without difficulty.

[^4]:    ${ }^{4}$ It is customary in such discussions to rule out a priori the possibility that roots may lie outside the unit circle; thus, the case $\left|z_{j}\right|>1$, the so called explosive case, is ruled out on a priori grounds.
    ${ }^{5}$ Since the matrix involved in $A(I)$ is $A(1)$, the condition that $A(I)$ is nonsingular rules out a real unit root; it does not necessarily rule out (pairs of) complex roots with unit modulus. This eventuality is generally ignored in practice, and we follow the prevailing custom.

