CIRCULANT MATRICES AND TIME-SERIES ANALYSIS

by D.S.G. Pollock

Queen Mary, University of London and

GREQAM: Groupement de Recherche en Economie Quantitative d'Aix-Marseille

Email: stephen pollock@sigmapi.u-net.com

This paper sets forth some of the salient results in the algebra of circulant matrices which can be used in time-series analysis. It provides easy derivations of some results that are central to the analysis of statistical periodograms and empirical spectral density functions. A statistical test for the stationarity or homogeneity of empirical processes is also presented.

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1. Introduction

A real circulant stochastic process of order T is one in which the autocovariance matrix of a vector $[x_0, x_1, \ldots, x_{T-1}]$, sampled from the process, is unchanged when the elements of the vector are subjected to a cyclical permutation. Thus, for any integers $s, t \in [0, T-1]$, the autocovariance should be

(1)
$$\gamma_{|t-s|} = C(x_t, x_s) = C\{x_{(t+\tau \mod T)}, x_{(s+\tau \mod T)}\},$$

where τ is an arbitrary integer. Such a process is the finite equivalent of a stationary stochastic process.

An ordinary stationary process is, by definition, distributed over the set of all positive and negative integers, which corresponds to a set of equally spaced points on the time axis. It is statistically invariant with respect to translations along this axis. A circular process, in comparison, is invariant with respect of translations around a closed circuit where the successive stations are indexed by the integers $0, 1, \ldots, T-1$.

One of the bugbears of time-series analysis is the need to cope with the fact that all data series are finite with a definite beginning and an end. Often, we need to infer the values of elements that lie outside the sample. Also, in attempting to determine the sampling properties of time-series estimates, we have often to contend with the disjunctions at the beginning and the end of the samples. Usually, we can show that, as the sample size increases, the effect of these disjunctions upon our estimates vanishes; but it is often difficult or laborious to demonstrate the fact.

A virtue of circular processes is that they suffer from none of these end effects. They are often easier to analyse than are the corresponding finite segments sampled from stationary processes, and the results that are obtained are usually perfect prototypes for the asymptotic results for stationary processes.

The theoretical importance of circular processes is reaffirmed once it is recognised that all time-series methods that make use of the discrete Fourier transform are effectively based upon the assumption that the data are from periodic processes, which are circular processes in other words.

The purpose of this paper is to set forth some of the salient results in the algebra of circulant matrices which can be used in time-series analysis. In the course of doing so, we shall provide easy proofs of some results that are central to the analysis of statistical periodograms and empirical spectral density functions. We shall also derive a statistical test for the stationarity or homogeneity of an empirical process. Once the algebra of circulant matrices has been digested, it should appear that the test is of a familiar classical nature.

Circulant matrices have represented a mathematical curiosity ever since their first appearance in the literature in a paper by Catalan [3]. The literature on circulant matrices, from their introduction until 1920, was summarised in four papers by Muir [7]–[10]. A recent treatise on the subject, which contains a useful bibliography, has been provided by Davis [5]; but his book does not deal with problems in time-series analysis.

Circulant matrices have been used, in passing, in the analysis of time series by several authors, notably by Anderson [1] and by Fuller [6], but their usefulness in organising the material of statistical Fourier analysis, or of time-series analysis in the frequency domain, has not been fully exploited. It is hoped that the paper will help in changing this.

2. Circulant Matrices and Polynomials

A circulant matrix is a Toeplitz matrix which has the general form of

(2)
$$X = \begin{bmatrix} x_0 & x_{T-1} & x_{T-2} & \dots & x_1 \\ x_1 & x_0 & x_{T-1} & \dots & x_2 \\ x_2 & x_1 & x_0 & \dots & x_3 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ x_{T-1} & x_{T-2} & x_{T-3} & \dots & x_0 \end{bmatrix}.$$

The vectors of such a matrix are generated by applying a succession of cyclic permutations to the leading vector, which therefore serves to specify the matrix completely. The elements of the circulant matrix $X = [x_{ij}]$, which may be real or complex numbers, fulfil the condition that $x_{ij} = x\{(i-j) \mod T\}$. Hence, the index for the supra-diagonal elements, for which $1 - T \leq (i-j) < 0$, becomes $(i-j) \mod T = T + (i-j)$.

The operator which effects the cyclic permutation of the elements of any (column) vector of order T is the matrix $K = [e_1, \ldots, e_{T-1}, e_0]$. This is formed from the identity matrix $I = [e_0, e_1, \ldots, e_{T-1}]$ by moving the leading vector to the back

of the array. Thus,

(3)
$$K = \begin{bmatrix} 0 & 0 & 0 & \dots & 0 & 1 \\ 1 & 0 & 0 & \dots & 0 & 0 \\ 0 & 1 & 0 & \dots & 0 & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\ 0 & 0 & 0 & \dots & 0 & 0 \\ 0 & 0 & 0 & \dots & 1 & 0 \end{bmatrix}.$$

We should observe that the following conditions hold:

(i)
$$K^{-q} = K^{T-q}$$
,
(ii) $K^0 = K^T = I$,
(iii) $K' = K^{T-1} = K^{-1}$.

Any circulant matrix X of order T can be expressed as a linear combination of the set of basis matrices I, K, \ldots, K^{T-1} ; and thus X can be expressed as a polynomial function

(5)
$$X = x_0 I + x_1 K + \dots + x_{T-1} K^{T-1} = x(K).$$

The conjugate of the matrix, obtained by replacing its elements by their complex conjugates, will be denoted by $\bar{X} = \bar{x}(K)$.

If $\{\gamma_i\}$ is an absolutely convergent sequence obeying the condition that $\sum |\gamma_i| < \infty$, then the z-transform of the sequence, which is defined by $\gamma(z) = \sum \gamma_j z^j$, is an analytic function on the unit circle. In that case, replacing z by K gives rise to a circulant matrix $\Gamma = \gamma(K)$ with finite-valued elements. Noting that $K \uparrow q = K \uparrow (q \mod T)$, it is found that

(6)
$$\Gamma = \left\{ \sum_{j=0}^{\infty} \gamma_{jT} \right\} I + \left\{ \sum_{j=0}^{\infty} \gamma_{(jT+1)} \right\} K + \dots + \left\{ \sum_{j=0}^{\infty} \gamma_{(jT+T-1)} \right\} K^{T-1}$$
$$= \psi_0 I + \psi_1 K + \dots + \psi_{T-1} K^{T-1}.$$

Given that $\{\gamma_i\}$ is a convergent sequence, it follows that the sequence of the matrix coefficients $\{\psi_0, \psi_1, \dots, \psi_{T-1}\}$ converges to $\{\gamma_0, \gamma_1, \dots, \gamma_{T-1}\}$ as T increases.

Notice that the matrix $\psi(K) = \psi_0 I + \psi_1 K + \cdots + \psi_{T-1} K^{T-1}$, which is derived from a polynomial $\psi(z)$ of degree T-1, is a synonym for the matrix $\gamma(K)$, which is derived from the z-transform of an infinite convergent sequence.

The polynomial representation is enough to establish that circulant matrices commute in multiplication and that their product is also a polynomial in K. That is to say

(7) If X = x(K) and Y = y(K) are circulant matrices, then their product W = XY = YX is also a circulant matrix.

The peculiar feature of such polynomials is that the powers of the argument K form a T-periodic sequence such that $K^{j+T} = K^j$; and thus the degrees of the polynomial products never exceed T-1. This periodicity is analogous to the periodicity of the powers of the argument $z = \exp\{-i2\pi/T\}$ which is to be found in the Fourier transform of a sequence of order T.

A symmetric circulant matrix G = G' must have the following representation:

(8)
$$G = c_0 I + c_1 (K + K^{-1}) + \dots + c_{T-1} (K^{T-1} + K^{1-T})$$
$$= c_0 I + (c_1 + c_{T-1}) K + \dots + (c_{T-1} + c_1) K^{T-1}.$$

The implication of the condition is that $g_{\tau} = (c_{\tau} + c_{T-\tau}) = g_{T-\tau}$.

Observe that this condition of matrix symmetry can be obtained directly from the defining conditions of a circulant matrix. For, if $x_{ij} = x\{(i-j) \bmod T\}$ is the element in the *i*th row and the *j*th column and if $x_{ij} = x_{ji}$, then $x\{\tau \bmod T\} = x\{-\tau \bmod T\}$ where $\tau = i-j$, which is to say that $x_{\tau} = x_{T-\tau}$.

Example 1. Imagine that T observations, running from t = 0 to t = T - 1, have been taken on a stationary and invertible ARMA(p,q) process y(t) which is described by the equation

$$(9) \qquad (1 + \alpha_1 L + \dots + \alpha_p L^p) y(t) = (1 + \mu_1 L + \dots + \mu_q L^q) \varepsilon(t),$$

wherein $\varepsilon(t)$ is a white-noise sequence of independently and identically distributed random variables of zero mean, and L is the lag operator which has the effect that $L\varepsilon(t) = \varepsilon(t-1)$. Corresponding to the observations, there is a set of T equations which can be arrayed in a matrix format:

$$\begin{bmatrix}
y_0 & y_{-1} & \dots & y_{-p} \\
y_1 & y_0 & \dots & y_{1-p} \\
\vdots & \vdots & \ddots & \vdots \\
y_p & y_{p-1} & \dots & y_0 \\
\vdots & \vdots & \ddots & \vdots \\
y_{T-1} y_{T-2} & \dots & y_{T-p-1}
\end{bmatrix}
\begin{bmatrix}
1 \\
\alpha_1 \\
\vdots \\
\alpha_p
\end{bmatrix} =
\begin{bmatrix}
\varepsilon_0 & \varepsilon_{-1} & \dots & \varepsilon_{-q} \\
\varepsilon_1 & \varepsilon_0 & \dots & \varepsilon_{1-q} \\
\vdots & \vdots & \ddots & \vdots \\
\varepsilon_q & \varepsilon_{q-1} & \dots & \varepsilon_0 \\
\vdots & \vdots & \ddots & \vdots \\
\varepsilon_{T-1} \varepsilon_{T-2} & \dots & \varepsilon_{T-q-1}
\end{bmatrix}
\begin{bmatrix}
1 \\
\mu_1 \\
\vdots \\
\mu_q
\end{bmatrix}.$$

Here, the generic equation is

(11)
$$\sum_{i=0}^{p} \alpha_i y_{t-i} = \sum_{i=0}^{q} \mu_i \varepsilon_{t-i}, \quad \text{where} \quad \alpha_0 = \mu_0 = 1.$$

Apart from the elements $y_0, y_1, \ldots, y_{T-1}$ and $\varepsilon_0, \varepsilon_1, \ldots, \varepsilon_{T-1}$ which fall within the sample period, these equations comprise the presample values y_{-p}, \ldots, y_{-1} and $\varepsilon_{-q}, \ldots, \varepsilon_{-1}$ which are to be found in the top-right corners of the matrices.

An alternative representation of the system of equations can be given which is in terms of polynomials. Thus, if

(12)
$$y(z) = y_{-p}z^{-p} + \dots + y_0 + y_1z + \dots + y_{T-1}z^{T-1}$$
$$\varepsilon(z) = \varepsilon_{-q}z^{-q} + \dots + \varepsilon_0 + \varepsilon_1z + \dots + \varepsilon_{T-1}z^{T-1},$$
$$\alpha(z) = 1 + \alpha_1z + \dots + \alpha_pz^p \quad \text{and}$$
$$\mu(z) = 1 + \mu_1z + \dots + \mu_qz^q,$$

then

(13)
$$y(z)\alpha(z) = \varepsilon(z)\mu(z).$$

By performing the polynomial multiplication of both sides of (13) and by equating the coefficients of the same powers of z, it will be found that the equation associated with z^t is precisely the generic equation under (11).

To derive the circulant representation of order T of the ARMA equations, we impose the conditions that $y_t = y_{(t \mod T)}$ and that $\varepsilon_t = \varepsilon_{(t \mod T)}$. In terms of equation (10), the effect of these conditions is to replace the presample elements by elements from within the sample. Thus, $y_{-1} = y_{T-1}, \ldots, y_{-p} = y_{T-p}$ and, likewise, $\varepsilon_{-1} = \varepsilon_{T-1}, \ldots, \varepsilon_{-q} = \varepsilon_{T-q}$.

With these provisos, the polynomials of equation (13) can be converted to circulant matrices simply by replacing the complex argument z by the matrix argument K.

3. Spectral Factorisation of Circulant Matrices

The matrix operator K has a spectral factorisation which is particularly useful in analysing the properties of the discrete Fourier transform. To demonstrate this factorisation, we must first define the so-called Fourier matrix. This is a symmetric matrix $U = T^{-1/2}[W^{jt}; t, j = 0, ..., T-1]$ whose generic element in the jth row and tth column is

(14)
$$W^{jt} = \exp(-i2\pi t j/T) = \cos(\omega_j t) - i\sin(\omega_j t),$$
 where $\omega_j = 2\pi j/T.$

On taking account of the T-periodicity of $W^q = \exp(-i2\pi q/T)$, the matrix can be written explicitly as

(15)
$$U = \frac{1}{\sqrt{T}} \begin{bmatrix} 1 & 1 & 1 & \dots & 1\\ 1 & W & W^2 & \dots & W^{T-1}\\ 1 & W^2 & W^4 & \dots & W^{T-2}\\ \vdots & \vdots & \vdots & & \vdots\\ 1 & W^{T-1} & W^{T-2} & \dots & W \end{bmatrix}.$$

The second row and the second column of this matrix contain the Tth roots of unity. The conjugate matrix is defined as $\bar{U} = T^{-1/2}[W^{-jt};t,j=0,\ldots,T-1];$ and, by using $W^{-q} = W^{T-q}$, this can be written explicitly as

(16)
$$\bar{U} = \frac{1}{\sqrt{T}} \begin{bmatrix} 1 & 1 & 1 & \dots & 1\\ 1 & W^{T-1} & W^{T-2} & \dots & W\\ 1 & W^{T-2} & W^{T-4} & \dots & W^2\\ \vdots & \vdots & \vdots & & \vdots\\ 1 & W & W^2 & \dots & W^{T-1} \end{bmatrix}.$$

It is readily confirmed that U is a unitary matrix fulfilling the condition

$$\bar{U}U = U\bar{U} = I.$$

This is proved using the formula for the partial sum of a geometric progression:

(18)
$$\sum_{t=0}^{T-1} W^{(j-k)t} = \begin{cases} \frac{1 - W^{T(j-k)}}{1 - W^{(j-k)}} = 0, & \text{if } j \neq k, \\ T, & \text{if } j = k. \end{cases}$$

The result follows, in both cases, from the fact that $W^{Tq}=W^0=1$.

It is helpful to note the following reciprocal relationships between the summation vector h = [1, 1, ..., 1]' and the vector $e_0 = [1, 0, ..., 0]'$ which is the leading column of the identity matrix:

(19)
$$(i) \quad h = \sqrt{T}Ue_0 = \sqrt{T}\bar{U}e_0,$$

$$(ii) \quad e_0 = \frac{1}{\sqrt{T}}\bar{U}h = \frac{1}{\sqrt{T}}Uh.$$

Now consider postmultiplying the unitary matrix U by the diagonal matrix

(20)
$$D = \begin{bmatrix} 1 & 0 & 0 & \dots & 0 \\ 0 & W^{T-1} & 0 & \dots & 0 \\ 0 & 0 & W^{T-2} & \dots & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & 0 & \dots & W \end{bmatrix}.$$

Then, it is easy to see, by the direct multiplication of the matrices of (15) and (20), that UD = KU, where K is the circulant operator from (3). From this, it follows that

(21)
$$(i) \quad K = UD\bar{U} = \bar{U}\bar{D}U,$$

$$(ii) \quad D = \bar{U}KU,$$

$$(iii) \quad \bar{D} = UK\bar{U},$$

where

(22)
$$\bar{D} = \text{diag}\{1, W, W^2, \dots, W^{T-1}\}\$$

is both the conjugate and the inverse of D. These equalities follow from the fact that K is a real-valued matrix which must equal its own conjugate.

Also, it can be shown that

(23)
$$K' = K^{-1} = U\bar{D}\bar{U} = \bar{U}DU.$$

The following conclusions can be reached in a straightforward manner:

(24) If
$$X = x(K)$$
 is a circulant matrix then

(i)
$$X = x(K) = Ux(D)\overline{U} = \overline{U}x(\overline{D})U$$
,

(ii)
$$X' = x(K') = Ux(\bar{D})\bar{U} = \bar{U}x(D)U$$
,

(iii)
$$X^{-1} = x^{-1}(K) = Ux^{-1}(D)\bar{U}.$$

4. Symmetric Circulant Matrices and their Fourier Transforms

We may describe the elements of the diagonal matrix $x(\bar{D})/T$ as the spectral ordinates of the circulant matrix X = x(K).

The jth diagonal element of the matrix $x(\bar{D})/T$ is given by

(25)
$$\xi_j = \frac{1}{T} \sum_t x_t W^{jt} = \frac{1}{T} \sum_t (x_t^{re} + ix_t^{im}) \left\{ \cos(\omega_j t) - i\sin(\omega_j t) \right\},$$

where $\omega_j = 2\pi j/T$ is described as the jth Fourier frequency. In fact, the sequence $\{\xi_0, \xi_1, \dots, \xi_{T-1}\}$ of the diagonal elements of $x(\bar{D})/T$ is nothing but the discrete Fourier transform (DFT) of the sequence $\{x_0, x_1, \dots, x_{T-1}\}$ of the elements in the leading vector of X = x(K).

As well as $x(\bar{D})$, we need to consider the matrix x(D) which features in the equations of (24) and the matrix $\bar{x}(D)$ which is the conjugate of $x(\bar{D})$.

(26) Let the jth element of $x(\bar{D})/T$ be denoted by ξ_j . Then the jth element of $\bar{x}(D)/T$ is ξ_j^* (i.e. the complex conjugate of ξ_j) and the jth element of x(D)/T is ξ_{T-j} .

These conditions can be used in establishing some basic results concerning real-valued circulant matrices and symmetric circulant matrices.

(27) If
$$X = \bar{X}$$
 is a real-valued circulant matrix such that $x(K) = \bar{x}(K)$, and $x(D) = \bar{x}(D)$, then, equivalently, $\xi_j^* = \xi_{T-j}$.

The implication is that, if $X = x(K) = \bar{x}(K)$ is a real-valued matrix, which is to say that $x_t^{im} = 0$ for all t, then the sequence of complex numbers $\{\xi_0, \xi_1, \dots, \xi_{T-1}\}$, defined according to (25), which are the diagonal elements of $x(\bar{D})/T$, consists of a real part which is an even or symmetric function of the index j and an imaginary part which is an odd or anti-symmetric function.

(28) If
$$X = X'$$
 is a symmetric circulant matrix such that $x(K) = x(K')$, and $x(D) = x(\bar{D})$, then, equivalently, $\xi_j = \xi_{T-j}$.

Here, the condition that $x(D) = x(\bar{D})$ arises out of the comparison of (24)(i) and (24)(ii).

If a circulant matrix is both real-valued and symmetric, then the elements of the corresponding Fourier transform will fulfil the same conditions as the elements of the matrix. That is to say

(29) If $X = x(K) = \bar{x}(K) = x(K')$ is a real-valued symmetric matrix with $x_{\tau} = x_{\tau}^* = x_{T-\tau}$, then, equivalently, the diagonal elements of $x(\bar{D})$ will constitute a real-valued even sequence with $\xi_j = \xi_j^* = \xi_{T-j}$.

Of course, the condition of matrix symmetry, which is that $x_{\tau} = x_{T-\tau}$, ensures that the terms in the sine function cancel within the sum defining ξ_j in (25). Thus, when X is a real symmetric circulant matrix, its spectral ordinates $\xi_j = \sum_t x_t \cos(\omega_j t)$ are the products of a cosine Fourier transform which therefore constitute a real-valued even sequence.

Example 2. Consider an ARMA process represented, in the form of an infiniteorder moving-average, by $y(t) = \phi(L)\varepsilon(t)$, where $\phi(L) = \mu(L)/\alpha(L)$ represents the series expansion of a rational function of the lag operator. The autocovariance generating function of the process is

(30)
$$\gamma(z) = \sigma^2 \phi(z) \phi(z^{-1}) = \sum_{\tau=0}^{\infty} \gamma_{\tau}(z^{\tau} + z^{-\tau}),$$

where $\gamma_{\tau} = \sigma^2 \sum_{j=0}^{\infty} \phi_j \phi_{j+\tau}$ is the autocovariance of lag τ . Setting $z = \exp\{-\mathrm{i}\omega_j\}$ and recognising that $\exp\{\mathrm{i}\omega_j\tau\} + \exp\{-\mathrm{i}\omega_j\tau\} =$ $2\cos(\omega_i\tau)$ gives an ordinate of the spectral density function in the form of

(31)
$$f(\omega_j) = \frac{1}{2\pi} \gamma(\exp\{i\omega_j\}) = \frac{1}{2\pi} \Big\{ \gamma_0 + 2\sum_{\tau=1}^{\infty} \gamma_\tau \cos(\omega_j \tau) \Big\},$$

which represents the cosine Fourier transform of the infinite sequence of ordinary autocovariances. The spectral density function is ordinarily regarded as a continuous function of the frequency variable $\omega \in [0, \pi]$. However, the spectral ordinates at the frequency values $\omega_j = 2\pi j/T$; $j = 0, \dots, T-1$, which are equally spaced across the interval $[0, 2\pi)$, constitute an even real-valued sequence.

The same spectral ordinates can be obtained from the elements of the diagonal matrix

(32)
$$\frac{1}{2\pi}\gamma(D) = \frac{1}{2\pi}\bar{U}\gamma(K)U = \frac{\sigma^2}{2\pi}\bar{U}\phi(K')\phi(K)U,$$

where $\gamma(K) = \Gamma$ and $\phi(K) = \Phi$ are the circulant matrices which are obtained when z is replaced by K in $\gamma(z)$ and $\phi(z)$ respectively. Observe that, whereas, in general, the coefficients of $\gamma(z)$ constitute an infinite sequence, the coefficients of $\gamma(D)$ form only a finite sequence of order T. The relationship between the coefficients of $\gamma(z)$ and those of $\gamma(D) = \psi(D)$ is expressed in equation (6).

5. The Fourier Transform and the Periodogram

Let $X = x(K) = Ux(D)\bar{U} = \bar{U}x(\bar{D})U$ be a real-valued circulant matrix. The DFT of the vector $x = x(K)e_0$, which is given by $T^{-1/2}Ux(K)e_0 = T^{-1/2}x(\bar{D})Ue_0 =$ $T^{-1}y(\bar{D})h$, can be written as

$$\xi = T^{-1/2}Ux,$$

Premultiplying the latter equation by $T^{1/2}\bar{U}$ gives

$$(34) x = T^{1/2}\bar{U}\xi;$$

which represents the inverse transform by which x is recovered from $\xi = x(\bar{D})h/T$. By referring to (25), it can be seen that the jth element of ξ is

(35)
$$\xi_{j} = \frac{1}{T} \sum_{t=0}^{T-1} x_{t} \exp(-i\omega_{j}t) = \frac{1}{T} \sum_{t=0}^{T-1} \left\{ x_{t} \cos(\omega_{j}t) - ix_{t} \sin(\omega_{j}t) \right\},$$

Observe that, with $\omega_j = 2\pi j/T$, there are $\cos(\omega_0 t) = 1$ and $\sin(\omega_0 t) = 0$ when j = 0. Also, if T is even, then T/2 is an integer and $\omega_{T/2} = \pi$. In that case, the sum comprises $\cos(\omega_{T/2} t) = (-1)^t$, whereas $\sin(\omega_{T/2} t) = 0$ vanishes.

Equation (35) is therefore amenable to various specialisations according to the value of j. First, when j = 0 the formula delivers the sample mean:

(36)
$$\xi_0 = \frac{1}{T} \sum_{t=0}^{T-1} x_t = \bar{x} = \alpha_0.$$

Next, in consequence of the real-valued nature of the data, there are

(37)
$$\xi_j = \frac{\alpha_j - i\beta_j}{2} \quad \text{and} \quad \xi_{T-j} = \xi_j^* = \frac{\alpha_j + i\beta_j}{2},$$

with

(38)
$$\alpha_j = \frac{2}{T} \sum_{t=0}^{T-1} x_t \cos(\omega_j t) \quad \text{and} \quad \beta_j = \frac{2}{T} \sum_{t=0}^{T-1} x_t \sin(\omega_j t),$$

where j = 1, ..., (T-1)/2 if T is odd and j = 1, ..., (T/2)-1 if T is even. Both cases are covered by setting j = 1, ..., [(T-1)/2], where $[(T-1)/2] = \text{trunc}\{(T-1)/2\}$ is the integer part of (T-1)/2. Finally, there is

(39)
$$\xi_{T/2} = \frac{1}{T} \sum_{t=0}^{T-1} (-1)^t x_t = \alpha_{T/2} \quad \text{if } T \text{ is even.}$$

The elements of the inverse Fourier transform can be written as

(40)
$$x_t = \sum_{j=0}^{T-1} \xi_j \exp(i\omega_j t) = \sum_{j=0}^{T-1} \left\{ \xi_j \cos(\omega_j t) + i\xi_j \sin(\omega_j t) \right\}.$$

However the conditions of (37) imply that

(41)
$$\alpha_j = \xi_j + \xi_j^* \quad \text{and} \quad \beta_j = -\mathrm{i}(\xi_j - \xi_j^*),$$

where j = 1, ..., [(T-1)/2]. Therefore, on the understanding that $\beta_{T/2} = 0$ if T is even, equation (40) can be written as

$$(42) x_t = \sum_{j=0}^{[T/2]} \left\{ \alpha_j \cos(\omega_j t) + \beta_j \sin(\omega_j t) \right\} = \alpha_0 + \sum_{j=1}^{[T/2]} \rho_j \cos(\omega_j t - \theta_j),$$

where $\rho_j^2 = \alpha_j^2 + \beta_j^2$ and $\theta_j = \tan^{-1}(\beta_j/\alpha_j)$, and where $[T/2] = \operatorname{trunc}(T/2)$. The product $G = T^{-1}X'X$ is the matrix of the so-called circular autocovari-

ances. Its elements $g_0, g_1, \ldots, g_{T-1}$ are given by the formula

$$g_{\tau} = \frac{1}{T} \sum_{t=0}^{T-1} x_t x_{t+\tau}; \quad \text{where} \quad x_t = x_{(t \mod T)}$$
or, equivalently,
$$g_{\tau} = \frac{1}{T} \sum_{t=0}^{T-1-\tau} x_t x_{t+\tau} + \frac{1}{T} \sum_{t=0}^{\tau-1} x_t x_{t+T-\tau}$$

$$= c_{\tau} + c_{T-\tau}.$$

If the mean \bar{x} of the data values $x_0, x_1, \ldots, x_{T-1}$ is non-zero, then its square is liable to be subtracted from the circular autocovariances. The circular autocovariances, which constitute an even sequence, differ from the ordinary autocovariances c_{τ} , which are obtained from the first formula by setting $x_t = 0$ when t > T - 1, instead of setting $x_t = x_{T-t}$, or from the second formula by taking only the leading sum. The circular autocovariances constitute an even sequence.

The matrix of circular autocovariances can be expressed variously as

(44)
$$G = \bar{U}g(D)U = T^{-1}X'X = T^{-1}\bar{U}x(\bar{D})x(D)U.$$

The core of this expression is the real-valued diagonal matrix

(45)
$$g(D) = T^{-1}x(\bar{D})x(D) = T\operatorname{diag}\{|\xi_0|^2, |\xi_1|^2, \dots, |\xi_{T-1}|^2\}.$$

Consider premultiplying equation (44) by e'_0 and postmultiplying it by e_0 . Since $e'_0Ge_0 = g_0$ and $Ue_0 = \bar{U}e_0 = T^{-1/2}h$, we get the following expression for the mean square of the data:

(46)
$$g_0 = T^{-2}h'x(\bar{D})x(D)h = \sum_{j=0}^{T-1} |\xi_j|^2.$$

The terms in the sum on RHS can be specialised in various ways. When j=0, there is

$$|\xi_0|^2 = \alpha_0^2 = \bar{x}^2.$$

When j = 1, ..., [(T - 1)/2], there are

(48)
$$|\xi_j|^2 = |\xi_{T-j}|^2 = \frac{\alpha_j^2 + \beta_j^2}{4} = \frac{\rho_j^2}{4}.$$

Finally, there is

(49)
$$|\xi_n|^2 = \alpha_n^2 \quad \text{if} \quad T = 2n \quad \text{is even.}$$

The terms $|\xi_j|^2$ and $|\xi_{T-j}|^2$ can be combined within the sum of equation (46). It follows that, when T is even, the equation can be rewritten as

(50)
$$g_0 - \alpha_0^2 = \sum_{j=0}^{n-1} \frac{\rho_j^2}{2} + \alpha_n^2.$$

When T is odd, the summation is extended to include the nth term which becomes $\rho_n^2/2 = (\alpha_n^2 + \beta_n^2)/2$. Observe that $g_0 - \alpha_0^2 = T^{-1} \sum_t (x_t - \bar{x})^2$ is the sample variance calculated from the mean-adjusted data.

The sequence $\{T\rho_1^2/2, T\rho_2^2/2, \ldots, T\rho_{n-1}^2/2, T\alpha_n^2\}$, which applies when T is even, together with the alternative sequence $\{T\rho_1^2/2, T\rho_2^2/2, \ldots, T\rho_n^2/2\}$, which applies when T is odd, constitute the periodogram of the data. Equation (50) represents an analysis of variance in terms of the periodogram.

Example 3. In equation (44), the matrix U represents a set of complex eigenvectors of the matrix G. The diagonal matrix g(D) contains the corresponding set of real-valued eigenvalues. A set of real-valued eigenvectors can be found to replace the complex ones. Given that G = g(K) and g(D) are a real-valued matrices, it follow from (44) that

(51)
$$GU = Ug(D)$$
 and $G\bar{U} = \bar{U}g(D)$.

By combining these equations in two ways, we can write

(52)
$$GC = Cg(D)$$
 and $GS = Sg(D)$,

where

(53)
$$C = \left(\frac{U + \bar{U}}{2}\right) \quad \text{and} \quad S = \left(\frac{U - \bar{U}}{2i}\right)$$

are, respectively, a matrix of cosines and a matrix of sines.

To reveal the structure of these matrices, we may adopt a notation whereby $c_p = [\cos(\omega_p t); t = 0, \dots, n]'$ and $s_q = [\sin(\omega_q t); t = 0, \dots, n]'$, where n = [T/2]. Also, let $h = [1, 1, \dots, 1]'$ and $j = [1, -1, \dots, (-1)^{T-1}]'$. Then

(54)
$$C = \begin{cases} [h, c_1, \dots, c_{n-1}, j, c_{n-1}, \dots, c_1], & \text{if } T \text{ is even;} \\ [h, c_1, \dots, c_{n-1}, c_n, c_n, \dots, c_1], & \text{if } T \text{ is odd.} \end{cases}$$

Likewise,

(55)
$$S = \begin{cases} [0, s_1, \dots, s_{n-1}, 0, -s_{n-1}, \dots, -s_1], & \text{if } T \text{ is even;} \\ [0, s_1, \dots, s_{n-1}, s_n, -s_n, \dots, -s_1], & \text{if } T \text{ is odd.} \end{cases}$$

Now define the matrix which comprises the distinct elements of C and S ordered by rising frequency. This is

(56)
$$V = \begin{cases} [h, c_1, s_1, \dots, c_{n-1}, s_{n-1}, j], & \text{if } T \text{ is even;} \\ [h, c_1, s_1, \dots, s_{n-1}, c_n, s_n], & \text{if } T \text{ is odd.} \end{cases}$$

The matrix V is square and of full rank, and it satisfies the equation

$$(57) GV = Vg(D).$$

Therefore, we can describe the vectors of sines and cosines within V as the characteristic vectors of the matrix of circular autocovariances.

When the characteristic vectors are ordered according to the decreasing sizes of their associated eigenvalues, they constitute the so-called principal components of the matrix X. The concept of principal components supplies additional meaning to the analysis of the periodogram.

An analysis of the principal components of the data has been the basis of a number of recent papers in paleoclimatology which purport to reveal the presence in meteorological phenomena of hidden cycles. (See, for example, Schlesinger and Ramankutty [12] and Vautard and Ghil [13].)

The procedure for uncovering the cycles, which is known as singular spectral analysis, is not based upon circulant matrices. Instead, it is based, typically, upon a data matrix which can be obtained from the leading matrix of equation (10) by replacing the presample elements by zeros and by appending an upper triangular matrix at the bottom whose nonzero elements are the remainder of the sample values.

The reality of the climatological cycles has been disputed; and those who claim to have uncovered them seem to be unable to provide convincing proof. There is a suspicion that the cycles that have been identified are merely artifacts of the data analysis. If the analysis were conducted using circulant matrices, then it should be straightforward to employ standard techniques to assess whether the periodogram ordinates are significantly different from those which might originate in a whitenoise process.

6. Complex Normal Distributions

In this section, we shall develop the statistical properties of a complex normal vector which is obtained by applying a Fourier transform to a spherical normal vector $x \sim N(0, \sigma^2 I)$, with a zero mean vector E(x) = 0 and with dispersion matrix $D(x) = \sigma^2 I$, which represents a segment of a white-noise process. The results will be used in the following section in analysing the statistical properties

of the periodogram of a normal vector $x \sim N(0, \sigma^2 Q)$ obtained by sampling an arbitrary stationary stochastic process.

The basic results, which are exact, are obtained under the assumption that x is from a circular process and that, accordingly, the dispersion matrix Q is a circulant matrix. These results provide asymptotically valid approximations for the more general and realistic cases where x is generated by a linear stochastic process which is not necessarily a normal process.

In order to develop the statistical theory of a complex normal vector, we need to generalise some of the familiar operations of multivariate statistical analysis.

First we define the expectation of a complex random vector:

(58) If ξ is a complex vector with a generic element of $\xi_j = (\alpha_j + i\beta_j)/2$, then its expectation $E(\xi)$ is a vector whose generic element is $E(\xi_j) = \{E(\alpha_j) + iE(\beta_j)\}/2$.

Next we define the dispersion matrix of a complex vector:

(59) If ξ is a complex vector of zero expectation with a generic element of $\xi_j = (\alpha_j + i\beta_j)/2$, then its variance-covariance matrix or dispersion matrix $D(\xi) = E(\xi \bar{\xi}') = Q$ has $E(\xi_j \xi_k^*) = E\{(\alpha_j + i\beta_j)(\alpha_k - i\beta_k)\}/4$ as the element in the j, kth position.

In general, the dispersion matrix of a complex vector is a Hermitian matrix $Q = \bar{Q}'$ which is equal to its own conjugate transpose. However, if $\xi = T^{-1/2}Ux$ is the Fourier transform of a real-valued vector, then the real and imaginary elements are uncorrelated, and its dispersion matrix Q = Q' is real and symmetric.

We are interested, particularly, in the statistical properties of the Fourier transform of a real-valued normal vector. The simplest case is that of a vector of independently and identically distributed elements:

(60) If $x \sim N(0, \sigma^2 I_T)$ is a real-valued vector with a spherical normal distribution and if U is a unitary transformation, then $T^{1/2}\xi = Ux \sim N_C(0, \sigma^2 I_T)$ has a complex spherical normal distribution.

Here the subscripted C on N_C indicates that this is the distribution of a complex vector.

Notice that the elements of ξ are not statistically independent in the ordinary sense since, if it is the transform of a real-valued vector x, then ξ will comprise both $\xi_j = (\alpha_j + i\beta_j)/2$ and its conjugate $\xi_{T-j} = (\alpha_j - i\beta_j)/2$. However, the complex coordinates $\alpha_j, \beta_j; j = 0, \ldots, n = [T/2]$ will be mutually uncorrelated, as will be the corresponding complex moduli $\rho_j^2 = \alpha_j^2 + \beta_j^2$.

It is straightforward to associate a chi-square distribution with the modulus of a complex normal vector with a spherical distribution.

(61) Let $T^{1/2}\xi = Ux \sim N_C(0, \sigma^2 I_T)$ be a complex spherical normal vector which represents the Fourier transform of a real-valued spherical normal vector $x \sim N(0, \sigma^2 I_T)$. Then $T\bar{\xi}'\xi/\sigma^2 = x'x/\sigma^2 \sim \chi^2(T)$.

According to Cochran's theorem [4], if $x \sim N(0, \sigma^2 I_T)$ and if $P = P' = P^2$ is a real-valued symmetric matrix of order T with rank(P) = r, then $x'Px/\sigma^2 \sim \chi^2(r)$

and $x'(I-P)x/\sigma^2 \sim \chi^2(T-r)$ are statistically independent chi-square variates. The following is a extension of the result which is appropriate to complex spherical normal vectors:

(62) Let $T^{1/2}\xi = Ux \sim N_C(0, \sigma^2 I_T)$ be a complex spherical normal vector obtained from the real-valued vector x via the unitary Fourier matrix U of (15). Let $S = S' = S^2$ be a real-valued selection matrix of order T with r units and T - r zeros on the diagonal, which form an even sequence, and with zeros elsewhere. Then $\bar{U}SU = P = P' = P^2$ is a real-valued symmetric idempotent circulant matrix with rank(P) = r; and it follows that $T\bar{\xi}'S\xi/\sigma^2 = x'\bar{U}SUx/\sigma^2 \sim \chi^2(r)$ and $T\bar{\xi}'(I - S)\xi/\sigma^2 = x'\bar{U}(I - S)Ux/\sigma^2 \sim \chi^2(T - r)$ are statistically independent.

Here, we should recall that, according to (29), the condition that the diagonal elements of S constitute a real-valued even sequence is sufficient to ensure that $P = \bar{U}SU$ is a real-valued matrix.

Example 4. Let $\xi = T^{-1/2}Ux$ be the Fourier transform of a real-valued vector $x \sim N(0, \sigma^2 I_T)$. Let $\xi_j = (\alpha_j + \mathrm{i}\beta_j)/2$ and let $S_j = e_j e_j' + e_{T-j} e_{T-j}'$, where $0 < j \le [(T-1)/2]$. Then S_j is a diagonal matrix of rank 2 with units on the diagonal in the positions j and T-j and with zeros elsewhere. The diagonal elements of S_j constitute a real and even sequence. Therefore, it follows from (29) that $\bar{U}S_jU$ is a real-valued symmetric idempotent circulant matrix. Hence

(63)
$$\frac{T}{\sigma^2} \bar{\xi}' S_j \xi = \frac{T}{\sigma^2} (\xi_j \xi_j^* + \xi_{T-j} \xi_{T-j}^*) \\ = \frac{T}{2\sigma^2} (\alpha_j^2 + \beta_j^2) = \frac{T}{2\sigma^2} \rho_j^2 \sim \chi(2).$$

If T=2n is even, then both $S_0=e_0e_0'$ and $S_n=e_ne_n'$ are matrices of rank 1 whose diagonal elements constitute even sequences. (An even sequence must have the same element in the positions j and T-j; and, with j=0 or j=n=T/2, the two positions coincide.) Therefore, both $\bar{U}S_0U$ and $\bar{U}S_nU$ are real-valued symmetric idempotent circulant matrices of rank 1, and it follows that

(64)
$$\frac{T}{\sigma^2}\xi_0\xi_0^* = T\frac{\alpha_0^2}{\sigma^2} = T\frac{\bar{x}^2}{\sigma^2} \sim \chi^2(1)$$
 and
$$\frac{T}{\sigma^2}\xi_n\xi_n^* = T\frac{\alpha_n^2}{\sigma^2} \sim \chi^2(1).$$

If T = 2n + 1 is odd, then $\beta_n \neq 0$ and there are $\xi_n = (\alpha_n + i\beta_n)/2$ and $\xi_{T-n} = (\alpha_n - i\beta_n)/2$, which are adjacent values within the vector ζ . These give rise to a $\chi^2(2)$ variate in the manner of (63).

7. Statistical Properties of the Periodogram

To see the significance for the analysis of the periodogram of the results of the foregoing section, let us consider the ARMA process of (9), which can be written

in the form of $y(t) = \phi(L)\varepsilon(t)$, where $\phi(L) = \mu(L)/\alpha(L)$ represents the series expansion of the rational operator. The process can also be represented, in terms of the z-transform polynomials of (12), by putting $y(z) = \phi(z)\varepsilon(z)$.

The analysis proceeds by using the circulant version of the ARMA process. By putting $\varepsilon_{(t \mod T)}$ in place of ε_t , which replaces the elements outside the sample by elements from within the sample, and by replacing the complex argument z by the matrix argument K, we derive the matrix equation $Y = \Phi \mathcal{E}$, where Y = y(K), $\Phi = \phi(K)$ and $\mathcal{E} = \varepsilon(K)$. The equation of the circulant process, together with that of its Fourier transform, can therefore be expressed as

(65)
$$y = \Phi \varepsilon = \mathcal{E} \phi \text{ and } T^{-1/2}Uy = \zeta_y = \phi(D)\zeta_{\varepsilon},$$

where $y = Ye_0$, $\phi = \Phi e_0$ and $\varepsilon = \mathcal{E}e_0$ are the leading vectors of their respective matrices.

We shall assume that $y \sim N(0, \sigma^2 Q)$, where $Q = \Phi' \Phi$, which implies that $\varepsilon = \Phi^{-1} y \sim N(0, \sigma^2 I)$. The results of the previous section can be applied directly to the Fourier transform of ε which is $\zeta_{\varepsilon} = \phi^{-1}(D)\zeta_{y}$.

Let $S_k = e_k e'_k$ be the selection matrix with a unit in the kth diagonal position and with zeros elsewhere. Then

(66)
$$\zeta_{\varepsilon k} \zeta_{\varepsilon k}^* = \bar{\zeta}_{\varepsilon}' S_k \zeta_{\varepsilon} \\ = \bar{\zeta}_{y}' \phi^{-1}(\bar{D}) S_k \phi^{-1}(D) \zeta_{y} \\ = \frac{\zeta_{yk} \zeta_{yk}^*}{q_k},$$

where q_k is the kth diagonal element of the matrix $\bar{U}QU = q(D) = \phi(\bar{D})\phi(D)$. Therefore, by putting $\zeta_{yj}\zeta_{yj}^*/q_j$ in place of $\xi_j\xi_j^*$ in equation (63), we get

(67)
$$\frac{T}{\sigma^2 q_j} \{ \zeta_{yj} \zeta_{yj}^* + \zeta_{y(T-j)} \zeta_{y(T-j)}^* \} = \frac{T}{2\sigma^2 q_j} \rho_{yj}^2 \sim \chi(2).$$

where j = 1, ..., [(T-1)/2]. If T = 2n is even, then, in addition,

(68)
$$\frac{T}{\sigma^2 q_0} \zeta_0 \zeta_0^* = T \frac{\alpha_0^2}{\sigma^2 q_0} = T \frac{\bar{y}^2}{\sigma^2 q_0} \sim \chi^2(1)$$
and
$$\frac{T}{\sigma^2 q_n} \zeta_n \zeta_n^* = T \frac{\alpha_n^2}{\sigma^2 q_n} \sim \chi^2(1).$$

If T = 2n + 1 is odd, then ζ_n and ζ_{T-n} are adjacent values within the vector ζ . These give rise to a $\chi^2(2)$ variate in the manner of (67).

Now recall that the expected value of a $\chi^2(r)$ variate is r whilst its variance is 2r. Also recall that, according to (31), the value of the spectral density function at the frequency ω_j is $f(\omega_j) = \gamma(\exp\{-i\omega_j\})/(2\pi) = \sigma^2 q_j/(2\pi)$. It follows that

(69)
$$E\left(\frac{T}{2}\rho_{yj}^2\right) = 2\pi f(\omega_j)$$
 and $V\left(\frac{T}{2}\rho_{yj}^2\right) = (4\pi)^2 f^2(\omega_j)$ if $0 < \omega_j < \pi$.

Also, at zero frequency, there is

(70)
$$E(T\bar{y}^2) = 2\pi f(0)$$
 and $V(T\bar{y}^2) = 2(2\pi)^2 f^2(0)$,

whereas, if T=2n is even, then, at the highest observable frequency, there is

(71)
$$E(T\alpha_n^2) = 2\pi f(\omega_n)$$
 and $V(T\alpha_n^2) = 2(2\pi)^2 f^2(\omega_n)$ where $\omega_n = \pi$.

There remains the question of the validity of approximating a linear ARMA process by a circular process. In the case of a white noise process, there is no effective distinction between a linear process and a circular process—the results for the periodogram analysis are identical.

In the case of an ordinary linear ARMA process, the generic element is $y_t = \sum_j \phi_j \varepsilon_{t-j}$, whereas, for a circular process, it is $y_t = \sum_j \phi_j \varepsilon_{(t-j \mod T)}$. The difference lies entirely in the replacement of the presample elements in the sequence $\{\varepsilon_t; t=0,\pm 1,\pm 2,\ldots\}$ by elements from within the sample.

It is clear that, if the sequence $\{\phi_j; j=0,1,2,\ldots\}$ of ARMA parameters is absolutely convergent, then the presample values of $\{\varepsilon_t\}$ will be of diminishing significance in the formation of y_t as they recede in time. On this basis, it is to be expected that, as the sample index t increases, the corresponding elements of the two processes will converge. Therefore, as the sample size T increases, the preponderance of the sample elements from the two process will converge, and this should lead to the convergence of their periodograms.

The formal proposition which expresses these notions is as follows:

(72) Let $y(t) = \phi(L)\varepsilon(t)$, where $\varepsilon(t)$ is a white-noise process such that $E(\varepsilon_t) = 0$, $V(\varepsilon_t) = \sigma^2$ for all t, and where the coefficients of the operator $\phi(L)$ are absolutely summable such that $\sum_{j=0}^{\infty} |\phi_j| < \infty$. Let $q_j = |\phi(z)|^2$ with $z = \exp\{-\mathrm{i}\omega_j\}$. Then the periodogram of y(t), based on T observations, can be represented, at the Fourier frequency ω_j , by

$$\frac{T}{2}\rho_{yj}^2 = \frac{T}{2}q_j\rho_{\varepsilon j}^2 + R_j,$$

where
$$E(R_j) \to 0$$
 and $V(R_j) \to 0$ as $T \to \infty$.

Proofs of the proposition can be found, *inter alia*, in the texts of Priestley [11], Brockwell and Davies [2] and Fuller [6]. However, these proofs make no explicit reference to the circular process.

It appears from (69) that, when they are scaled by a factor of $1/(2\pi)$, the ordinates of the periodogram of a circular process are unbiased estimates of the corresponding ordinates of the spectral density function. It also appears that the variances of these estimates do not diminish as the sample size T increases. The same is roughly true of the ordinates of the periodogram of a linear process for which the results of (69) are asymptotically valid. The upshot is that the raw periodogram is not a consistent estimator of the spectral density function.

Nevertheless, as the sample size increases, the points at which the periodogram is defined become increasingly dense in every fixed frequency interval. Therefore,

one can generate ever-improving estimates of the underlying spectral density function by averaging an increasing number of adjacent periodogram ordinates from within an ever-decreasing frequency interval.

Let $M = U\mu(D)\bar{U}$ be a circulant smoothing matrix which serves as an averaging operator in respect of the periodogram, and let g(D) be the diagonal matrix of (45) whose elements give rise to the periodogram ordinates. Then s = Mg(D)h is a vector whose elements are the ordinates of an estimated spectral density function based on local averages of the periodogram.

An alternative expression for the vector of spectral estimates is $s = T^{1/2}U\mu(D)Ge_0$. This indicates that the estimated spectrum may be obtained by applying the relevant weights, contained within the diagonal matrix $\mu(D)$, to the vector $g = Ge_0$ of the circular autocovariances and then applying a Fourier transform to the product.

The equivalence of the two expressions is demonstrated by confirming the following identities, which make use of (19)(i):

(73)
$$Mg(D)h = U\mu(D)\bar{U}g(D)h$$
$$= U\mu(D)\bar{U}UG\bar{U}h$$
$$= T^{1/2}U\mu(D)Ge_0.$$

8. Testing the Homogeneity of a Time Series

The purpose of this final section is to show how the results of the foregoing sections can be used in a simple way to construct statistical tests of hypotheses relating to the periodogram, which are of a strictly classical nature.

Our specific object is to determine whether two adjacent segments of a time series have been generated by the same invariant process. Under the null hypothesis that the process is invariant, we may estimate a time-series model which can be used in constructing a filter for the purpose of reducing the series to white noise. Thereafter, we can test the null hypothesis that the two adjacent segments are generated by the same white-noise process.

An alternative to filtering the data is to divide the ordinates of the periodogram of the data by (estimates of) the corresponding spectral ordinates, in the manner of (66). The estimates of the spectral ordinates may be based upon the parameters of a fitted ARMA model. In that case, we should be taking the route to the spectrum which has been pursued in the Example 2. Alternatively, the spectral ordinates may be estimated by taking local averages of the periodogram ordinates in the manner indicated at the end of the previous section.

We shall avoid the unnecessary encumbrance of a notation which indicates explicitly that the vectors y_1 and y_2 , which are from a putative white-noise process, have been obtained by filtering the data or by weighting the periodogram ordinates. Therefore, according to the null hypothesis, $y_1 \sim N(0, \sigma^2 I_N)$ and $y_2 \sim N(0, \sigma^2 I_N)$ are two adjacent segments both comprising N elements, which are generated by a white-noise process which has a variance of σ^2 . Then the corresponding Fourier

transforms can be calculated for both segments independently via a scheme which is represented by the equation

$$\begin{bmatrix} \zeta_1 \\ \zeta_2 \end{bmatrix} = \begin{bmatrix} U & 0 \\ 0 & U \end{bmatrix} \begin{bmatrix} y_1 \\ y_2 \end{bmatrix},$$

where, compared with (33), a factor of $N^{1/2}$ has been incorporated in ζ_1 and ζ_2 in order to simplify the notation. According to the hypothesis, $\bar{\zeta}_1'\zeta_1/\sigma^2 \sim \chi^2(N)$ and $\bar{\zeta}_2'\zeta_2/\sigma^2 \sim \chi^2(N)$ are mutually independent chi-square variates and their sum $(\bar{\zeta}_1'\zeta_1 + \bar{\zeta}_2'\zeta_2)/\sigma^2 \sim \chi^2(2N)$ is also a chi-square variate.

Now consider forming the vector

(75)
$$\xi = \frac{1}{\sqrt{2}} \begin{bmatrix} U & U \end{bmatrix} \begin{bmatrix} y_1 \\ y_2 \end{bmatrix}.$$

We can show easily that $\bar{\xi}'\xi/\sigma^2 \sim \chi^2(N)$ is a chi-square variate. Moreover, by invoking Cochran's theorem, it can be show that, under the null hypothesis,

(76)
$$\frac{\bar{\xi}'\xi}{\sigma^2} \sim \chi^2(N) \quad \text{and} \quad \frac{\bar{\zeta}'_1\zeta_1 + \bar{\zeta}'_2\zeta_2 - \bar{\xi}'\xi}{\sigma^2} \sim \chi^2(N)$$

are mutually independent chi-square variates. Their ratio is therefore an F(N, N) variate; and this statistic provides the means of testing the hypothesis that y_1 and y_2 have been generated by the same process.

The relevant version of Cochran's theorem is as follows:

(77) Let $y \sim N(0, \sigma^2 I_{2N})$, and let $P_1 + P_2 = I_{2N}$ be a sum of two matrices which are symmetric and idempotent with $\operatorname{rank}(P_1) = \operatorname{rank}(P_2) = N$ such that $P_i = P_i^2$ and $P_i P_j = 0$ when $i \neq j$. Then $y' P_i y \sim \chi^2(N)$; i = 1, 2 are independent chi-square variates such that $(y' P_1 y + y' P_2 y)/\sigma^2 = y' y/\sigma^2 \sim \chi^2(2N)$.

To confirm that these conditions are fulfilled, we note that U is a unitary matrix such that $\bar{U}U = I$, where \bar{U} is the conjugate of U, and we note that

(78)
$$P_1 = \frac{1}{2} \begin{bmatrix} \bar{U} \\ \bar{U} \end{bmatrix} \begin{bmatrix} U & U \end{bmatrix} = \frac{1}{2} \begin{bmatrix} I & I \\ I & I \end{bmatrix}$$

is a symmetric idempotent matrix. We define $P_2 = I - P_1$ and we note that $P_1P_2 = 0$. It is easy to see in reference to (74) that

(79)
$$\frac{1}{\sigma^2} (\bar{\zeta}_1' \zeta_1 + \bar{\zeta}_2' \zeta_2) = \frac{1}{\sigma^2} (y_1' \bar{U} U y_1 + y_2' \bar{U} U y_2) \\ = \frac{1}{\sigma^2} y' y \sim \chi^2(2N).$$

Now we see, in view of the properties of P_1 and P_2 , that

(80)
$$\frac{\bar{\xi}'\xi}{\sigma^2} = \frac{y'P_1y}{\sigma^2} \quad \text{and} \quad \frac{\bar{\zeta}'_1\zeta_1 + \bar{\zeta}'_2\zeta_2 - \bar{\xi}'\xi}{\sigma^2} = \frac{y'P_2y}{\sigma^2}$$

are mutually independent chi-square variates of N degrees of freedom each.

The F statistic, which is formed from the ratio of the chi-square variates of (76), provides an overall test of the hypothesis of homogeneity. It will serve to detect cases where the power or variance differs between the two data segments. However, the statistic will not detect changes in the spectral structure of the process that do not affect the overall power.

Changes in the distribution of the power can be detected by making comparisons of the periodograms of the adjacent data segments in a number of disjoint frequency bands. The relevant statistics may be obtained from the existing F statistic by breaking the sum in the numerator into parts which constitute mutually independent chi-square variates.

Let S be a selection matrix of order $N \times N$ with zeros and units on the principal diagonal and with zeros elsewhere. Let the matrix be subject to the condition that, if it contains a unit in the jth position, then it also has a unit in the position N-j. If, in addition, the units are contiguous within two separate blocks, then the matrix will serve to select from the vectors ζ_1 , ζ_2 and ξ all of the ordinates which fall in a specific frequency band.

By applying the selection matrix to equations (74), we get

$$\begin{bmatrix} S\zeta_1 \\ S\zeta_2 \end{bmatrix} = \begin{bmatrix} SU & 0 \\ 0 & SU \end{bmatrix} \begin{bmatrix} y_1 \\ y_2 \end{bmatrix}.$$

By applying the selection matrix to equations (75), we get

(82)
$$S\xi = \frac{1}{\sqrt{2}} \begin{bmatrix} SU & SU \end{bmatrix} \begin{bmatrix} y_1 \\ y_2 \end{bmatrix}.$$

Let $S\zeta_1 = \zeta_{S1}$, $S\zeta_2 = \zeta_{S2}$ and $S\xi = \xi_S$. Then, under the hypothesis that there are no changes in the spectral density function within the selected range of frequencies, we should have

(83)
$$\frac{\bar{\zeta}'_{S1}\zeta_{S1} + \bar{\zeta}'_{S2}\zeta_{S2} - \bar{\xi}'_{S}\xi_{S}}{\bar{\xi}'\xi} \sim F(q, N), \quad \text{where} \quad q = \text{rank}(S).$$

To establish this result, it is enough to show that the numerator and the denominator of the statistic are quadratic forms in the same spherical normal vector which are based on mutually orthogonal symmetric idempotent matrices with ranks that correspond to the degrees of freedom of the F distribution.

The quadratic form of the denominator, which can be expressed as $y'P_1y$ where $y \sim N(0, \sigma^2 I_N)$, is based on the symmetric idempotent matrix $P_1 = P_1' = P_1^2$ of (78) which has rank $(P_1) = N$. The quadratic form of the numerator, which can be expressed as $y'P_Sy$, is based on the matrix

(84)
$$P_S = \frac{1}{2} \begin{bmatrix} \bar{U}SU & -\bar{U}SU \\ -\bar{U}SU & \bar{U}SU \end{bmatrix}.$$

It can be seen that $P_S = P'_S = P_S^2$ and that $P_1 P_S = 0$. Moreover, rank $(P_S) = \text{rank}(S) = q$. Thus, the conditions in question are fulfilled.

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