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By integrating out variables, $\left\{U_{i}\right\}$ is a set of mutually independent standardized normal variablea and so is $\left\{W_{i}\right\}$. Any square summable function $g\left(U_{1}, U_{2}, \ldots, U_{p}\right)$ can be approximated arbitrarily closely by the product set of orthonormal functions on the joint distribution of the $U_{i}$. This is true on the combined distribution of the $U_{i}$ and $W_{f}$ since both $g\left(u_{1}, u_{1}, \ldots, u_{p}\right)$ and the approximating series are constant for a given set $\left\{u_{1}, u_{n}, \ldots, u_{p}\right\}$.

Let us write the members of the product set on the distribution of the $U$ 's as $P_{i}$. We might take $P_{0}, P_{1}, P_{k}, \ldots$ to be, first a constant term, then the $u_{k}^{(1)}$, then the $u_{k}^{(2)}$, then the products $u_{k}^{(1)} u_{k^{\prime}}^{(2)} \ldots$. The $Q_{0}, Q_{1}, Q_{2}, \ldots$ may be similarly defined on the distribution of the $W$ 's. We now have

$$
\left.\begin{array}{lll}
\xi_{1}=\Sigma a_{i} P_{i}, & a_{0}=0, & \Sigma a_{i}^{2}=1 \\
\eta_{1}=\Sigma b_{i} Q_{j}, & b_{0}=0, & \Sigma b_{j}^{2}=1,
\end{array}\right\}
$$

and have to maximize $E\left(\xi_{1} \eta_{1}\right)$. Note that

$$
E\left\{\prod_{i=1}^{p} u_{i}^{(j)} \prod_{i=1}^{q} w_{i}^{(k)}\right\}=\prod_{i=1}^{p} E\left\{u_{i}^{(j)} w_{i}^{(k)}\right\}_{i-p+1} \prod_{i}^{q} E w_{i}^{(u)}
$$

since $\quad\left\{U_{1}, W_{1}\right\},\left\{U_{\mathbf{1}}, W_{2}\right\}, \ldots,\left\{U_{p}, W_{p}\right\}, W_{p+1}, \ldots, W_{q}$ are jointly independent.

But $E\left(u_{i}^{(H)} w_{i}^{(k)}\right)=\rho^{\prime} \delta_{j k}$ and so the only terms of (3•17) which do not vanish yield correlations of the form $\rho_{1}^{\left(i_{1}\right)} \rho_{\mathbf{g}}^{\left(\mathbf{v}_{1}\right)} \ldots \rho_{p}^{\left(\boldsymbol{L p}_{p}\right)}$. But of all such products $\rho_{1}$ is the greatest. If $\rho_{k}>\rho_{1}^{2}, \rho_{2}, \rho_{3}, \ldots, \rho_{k}$ are the next in order of magnitude.

By Lemma 1, the first pair of canonical variables is $U_{1} \equiv u_{1}^{(1)}, W_{1} \equiv w_{1}^{(1)}$, and the canonical correlation is $\rho_{1}$. Thus $\xi_{1}$ and $\eta_{1}$ must not contain any term in $u_{1}^{(1)}$ or $w_{1}^{(1)}$ respectively.
If $\rho_{\mathbf{1}}>\rho_{1}^{1}$ then the second pair of canonical variables is $U_{2} \equiv u_{2}^{(1)}, W_{\mathbf{2}} \equiv w_{3}^{(1)}$. The process can be continued as long as $\rho_{k}>\rho_{1}^{2}$. If $\rho_{p}>\rho_{1}^{2}$, then there are $p$ canonical variables which are linear in the $U_{i}$ or $W_{i}$ and hence in the $X_{i}$ and $Y_{i}$ respectively. In any case, at the $(p+1)$ th or earlier choice of pair, it will be necessary to select $u_{i}^{(s)}$ and $w_{1}^{(s)}$ and the functions are no longer linear in the $X_{i}$ and $Y_{i}$.

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## A $Q$-technique for the calculation of canonical variates

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## Summary

A $Q$-technique for evaluating canonical variates is described. It is shown to have computational, and statistical, advantages over the usual $R$-technique.

The canonical variates associated with $k$ multivariate normal populations, based on $v$ variates with common dispersion matrix $\Omega$ and means $\Gamma(k \times v)$, are linear combinations of the original variates. The coefficients $\xi_{i}$ of the $i$ th canonical variate are the elements of the latent vector associated with the ith largest latent root $\lambda_{i}$ of the equation $\left|\Gamma^{\prime} \Gamma-\lambda \Omega\right|=0$, where the means are supposed measured from an
origin such that each column of $\Gamma$ has zero sum. There are thus $v$ sets of these coefficients and they can be taken as columns of a square matrix $\Xi=\left(\xi_{1}, \xi_{2}, \ldots, \xi_{0}\right)$. The latent root and vector relationship can be written

$$
\begin{equation*}
\Gamma^{\prime} T \Xi=\Omega \Xi \Lambda_{1} \tag{I}
\end{equation*}
$$

where $\Lambda_{1}$ is the diagonal matrix of the latent roots.
The means $I$ II of the canonical variates are given by

$$
\begin{equation*}
\Pi=\Gamma \Xi \tag{2}
\end{equation*}
$$

The rows of $I I$ can be regarded as the co-ordinatee of points representing the means, referred to orthogonal canonical variate axes, and the square of the ordinary Euclidean distance between two means $\gamma_{i}$ and $\gamma_{j}$ is $\left(\gamma_{i}-\gamma_{j}\right) \boldsymbol{\Xi} \Xi^{\prime}\left(\gamma_{i}-\gamma_{j}\right)^{\prime}$. Thus if the vectors are acaled so that $\boldsymbol{\Xi} \Xi^{\prime}=\boldsymbol{\Omega}^{-1}$, i.e. $\boldsymbol{\Xi}^{\prime} \boldsymbol{\Omega} \boldsymbol{\Xi}=\mathbf{I}$, this distance becomes Mahalanobis's $\Delta_{i j}^{2}$. The axes corresponding to the $r$ largest latent roots maximize the total $\Delta^{\mathbf{2}}$ in $r$ dimensions between all pairs of populations (Reo, 1952, p. 365) or, equivalently, the first $r$ axes are the principal components of a set of $k$ pointe, representing the means, whose $\frac{1}{2} k(k-1)$ squared Euclidean distances are the values of $\Delta^{2}$ (Gower, 1986).
Alternatively, II can be found directly as the matrix of latent vectors of the $k \times k$ symmetric matrix $\boldsymbol{\theta}=\Gamma \boldsymbol{\Omega}^{-1} \Gamma^{\prime}$ provided that the vectors are scaled to make $\Pi^{\prime} \boldsymbol{I}=\Lambda_{\mathbf{1}}$, the diagonal matrix of latent roots of $\Theta$ (Gower, 1966). Thus

$$
\begin{equation*}
\theta I I=\Pi \Omega_{\mathbf{2}} . \tag{3}
\end{equation*}
$$

The matrices $\Lambda_{1}$ and $\Lambda_{1}$ will have the same non-zero diagonal elementa, there are $\min (k-1, v)$ of these, and will have zeros appended as necessary so that $\Lambda_{1}$ is $v \times v$ and $\Lambda_{1}$ is $k \times k$. The matrix $\Xi$ can be expressed in terms of $\Gamma$ and II by pre-multiplying both sides of (2) by $\Gamma^{\prime}$ and using (1) to give

Therefore

$$
\begin{align*}
\Gamma^{\prime} \boldsymbol{\Pi} & =\Gamma^{\prime} \Gamma \Xi=\Omega \Xi \Lambda_{1} \\
\Xi & =\Omega^{-1} \Gamma^{\prime} \Pi \Omega_{1}^{-} \tag{4}
\end{align*}
$$

where $\Lambda_{1}^{-}$is the generalized inverse of $\Lambda_{1}$ found by replacing all non-zero valuee by their reciprocals. In the special case when $\Gamma^{\prime} \Gamma$ is non-singular, i.e. $v<k$,

$$
\begin{equation*}
\Xi=\left(\Gamma^{\prime} \Gamma^{-1} \Gamma^{\prime} \Pi\right. \tag{5}
\end{equation*}
$$

When calculating canonical variates, sample values must replace the population parameters. The pooled within population dispersion matrix $W$ replaces $\boldsymbol{\Omega}$ and the sample means $G$ replace $\Gamma$. The usual method of computation is the $R$-technique derived from (1) and (2). It is then necessary to compute the vectors of ( $\left.G^{\prime} \mathbf{G}-\lambda \mathbf{W}\right) \mathbf{x}=0$ which requires either the vectors of the $v \times v$ non-symmetric matrix $\mathbf{W}^{-1} \mathbf{G}^{\prime} \mathbf{G}$ or a special calculation to factorize $\mathbf{W}$ into components $\mathrm{UU}^{\prime}$, where $\mathbf{U}$ is an upper triangular matrix (Ashton, Healy \& Lipton, 1957). Both of these methods of computation are more complex than the $Q$-technique derived from (3) and (4) which needs only the roots and vectors of the $k \times k$ symmetric matrix $T=G^{-1} \mathbf{G}^{\prime}$ and the standard matrix operations of inversion and multiplication. When $k<v$ the $Q$-teahnique has the additional advantage of requiring the vectors of a amaller matrix than the $R$-technique. The equation corresponding to (3) is

$$
\begin{equation*}
\mathbf{T P}=\mathbf{P L}_{\mathbf{2}} \tag{6}
\end{equation*}
$$

whose solution gives $\mathbf{P}$ and $L_{2}$ and thus also $L_{1}$. The canonical variate coefficients $\mathbf{X}$ are then given by equation (7), which corresponds to (4)

$$
\begin{align*}
& \mathbf{X}=\mathbf{W}^{-1} \mathbf{G}^{\prime} \mathbf{P} \mathbf{L}_{1}^{-}  \tag{7}\\
& \mathbf{X}=\left(\mathbf{G}^{\prime} \mathbf{G}\right)^{-1} \mathbf{G}^{\prime} \mathbf{P} \tag{8}
\end{align*}
$$

or when $\mathbf{G}^{\prime} \mathbf{G}$ is non-singular
Other advantages of the $Q$-technique are that the formula $D_{i j}^{2}=t_{i i}+t_{j j}-2 t_{i j}$ gives a simple direct method of computing the generalized distances from the elements of $T$, and it is possible to take account of the bias in $D^{2}$ (Rao, 1952, p. 364). Thus a matrix $\mathbf{E}$ of $D^{\mathbf{2}}$ values adjusted for bias by Rao's formula can be constructed. From this we form a new matrix $F$ with elements $-\frac{1}{2} e_{i j}$ and finally define $T_{1}$ as a matrix with elementa $t_{i j}^{11}=f_{i j}-f_{i}-f_{. j}+f .$. (Gower, 1866). This matrix $\mathbf{T}_{1}$ can be analysed as $\mathbf{T}$ above.
Canonical variates are often used as discriminant functions derived by finding the linear functions that give stationary values to the ratio of between-population to within-population sums of squares. This requires the solution of

$$
\begin{equation*}
(\mathbf{B}-\nu \mathbf{W}) \mathbf{x}_{1}=0 ; \tag{9}
\end{equation*}
$$

here $B$ is the usual sum of squares and products matrix between groups obtained by weighting each group mean by the sample size, whereas $\mathbf{G}^{\prime} \mathbf{G}$ is the corresponding unweighted sum of squares and
products matrix. If canonical variates are derived from (9) to give a set of coefficients $\mathbf{X}_{1}$ then scading the vectors so that $\mathbf{X}_{\mathbf{1}} \mathbf{X}_{1}^{\prime}=W^{-1}$ ensures that the squares of the distances between the means of the canonical variates are still equal to $D_{i s}^{2}$, but it is no longer generally true that the first $r$ dimensions maximize the total $D^{\mathbf{2}}$ in $r$ dimensions. The degree of approximation depends on the difference in the sample sizes; when all samples are the same size the two methods give the same results. These observations are relevant when canonical variates are used mainly for descriptive purposes (see, for example, Ashton, Healy \& Lipton, 1957), when it is reasonable to require a representation of the sample means which maximizes total $D^{2}$, using only a few dimensions.

The distribution functions of the diagonal elements of $\mathbf{B}$ are each proportional to $\chi^{*}$, but this is not so for $G^{\prime} G$, so that Bartlett's (1938) approximate test for the significance of canonical variates has an extra degree of approximation when $\mathbf{G}^{\prime} \mathbf{G}$ is used. Thus when $\mathbf{B}$ is used the $D^{\mathbf{2}}$ interpretation becomes approximate but when $\mathbf{G}^{\prime} \mathbf{G}$ is used the significance test is less reliable.

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## The non-null distribution of a statistic in principal components analysis

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## Sumbary

The non-null distribation is obtained of a teat statistic for the null hypothesis that a matrix with a single non-isotropic principal component has that component in a prescribed direction.

## 1. Introduotion

Suppose that one latent root of $\Sigma$, the dispersion matrix of $p$ multinormal variables with zero means

$$
\mathbf{x}^{\prime}=\left[x_{1}, \ldots, x_{p}\right]
$$

is $\sigma_{1}^{2}$ and the remaining roots $\sigma_{9}^{2}, \ldots, \sigma_{p}^{2}$ are all equal to $\sigma^{2}<\sigma_{1}^{2}$. If the latent column vector corresponding to $\sigma_{1}^{2}$ is $\mathbf{h}_{1}$, the variable $h_{1}^{\prime} \mathbf{x}$ is called the single non-isotropic principal component of $\Sigma$, because the variance of any linear combination of the variables $x$, orthogonal to $h_{1}^{\prime} \mathbf{x}$, is the same, namely $\sigma^{\mathbf{2}}$. Practical cases, where one comes across such a situation, have already been described (Kshirsagar, 1961). Given that $\Sigma$ admits such a single non-isotropic principal component and given a hypothetical direction vector $h$, it is of interest to test whether the hypothetical function $h^{\prime} x$ is the true non-isotropic principal component. In other words, one wishes to test the hypothesis

$$
H_{0}: \mathbf{h} \equiv \mathbf{h}_{\mathbf{1}} .
$$

The appropriate statistic for this is (Kshirsagar, 1981)

$$
\chi_{d}^{\mathbf{2}}=\frac{\mathbf{h}^{\prime} A^{t h}}{\mathbf{h}^{\prime} A \mathbf{h}}-\mathbf{h}^{\prime} A h,
$$

where

$$
A=X^{\prime} X
$$

and

$$
X=\left[x_{t r}\right] \quad(i=1, \ldots, p ; r=1, \ldots, n)
$$

