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A DISCRIMINANT RULE UNDER TRANSFORMATION

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Abstract We present a new rule for discriminating among continuous populations which are no multivariate normal. The basic idea is to construct the sample maximum likelihood discriminant rule after transforming the data by a suitable multivariate transformation to normality								
Key words: Cross-validation; Multivariate Box-Cox Transformation.								

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1. INTRODUCTION

Discriminant analysis is a widely applied technique in multivariate data analysis. Given g groups, populations or classes, G_1 , G_2 , ..., G_g , with corresponding probability density functions $f_1(x)$ on \mathbb{R}^p , i=1, ..., g, the problem is allocating an individual to one of these groups on the basis of his measurements $\mathbf{x}=(\mathbf{x}_1, \ldots, \mathbf{x}_p)$, which are considered as an observation from a population described by a random vector $\mathbf{X}=(X_1, X_2, \ldots, X_p)$, such that $\mathbf{X}\sim f_1$ under G_1 , i=1, ..., g. The allocation should be "optimal" in the sense of minimizing, on average, the number of incorrect assignments. When the densities $f_1(x)$ are known, a suitable discriminating rule is the maximum likelihood (ML) discriminant rule (Mardia et al. (1979), p. 301), defined as follows: Assign \mathbf{x} to G_1 if

$$f_{i}(\mathbf{x}) = \max_{1 \le j \le g} f_{j}(\mathbf{x}). \tag{1}$$

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We will assume, for the rest of this work, that the a priori probabilities of the classes are equal. When this is the case, it can be that the rule (1) minimizes the total probability shown misclassification (see Seber (1984), p. 331). When the densities $f_i(x)$ depend on unknown parameters θ_i , i.e., if $f_i(x) = f_{\theta_i}(x)$, i=1, ...,g, we use the sample ML (SML) discriminant rule (Mardia et al. (1979), p. 309), which is obtained by replacing, in rule (1), the parameters θ_i by efficient estimates $\theta_i = \theta_i(x_i)$, where x_i is a $n_i \times p$ data matrix from G_i , i=1, ...,g.

An important example of the ideas above appears when $f_i \sim N_p(\mu_i, \Sigma_i)$, and the parameters (μ_i, Σ_i) are unknown, i=1, ...,g. Let \overline{x}_i and S_i be, respectively, the sample mean and the covariance matrix of \mathfrak{X}_i , i=1, ...,g. The SML rule can adopt two possible different forms: (i) If $\Sigma_1 = \Sigma_2 = \ldots = \Sigma_g = \Sigma$, say, we assign x to G_i if

$$\ell_{i}^{\prime}\mathbf{x}+c = \max_{1 \leq i \leq g} [\ell_{j}^{\prime}\mathbf{x}+c_{j}^{\prime}], \tag{2}$$

where $\ell_1 = S_p^{-1} \overline{x}_1$, $c_1 = -(1/2) \overline{x}_1^2 S_p^{-1} \overline{x}_1$, i = 1, ..., g, and S_p is the pooled estimator $S_p = \sum_{l=1}^g (n_l - 1) S_l / (N - g)$, $N = \sum_{l=1}^g n_l$; (ii) If the dispersion matrices are unequal, we assign x to G_1 if

$$Q_{i}(\mathbf{x}) = \min_{1 \le j \le g} Q_{j}(\mathbf{x}), \tag{3}$$

where $Q_1(\mathbf{x}) = (1/2)\log[|S_1|] + (1/2)[(\mathbf{x} - \overline{\mathbf{x}}_1)'S_1^{-1}(\mathbf{x} - \overline{\mathbf{x}}_1)], \quad i=1, \dots, g.$

It is well-known that, when g=2, the rule (2) is equivalent to the Fisher's Linear Discriminant Function obtained by Fisher (1936) not assuming normality but looking for a "sensible" rule for discriminating between the groups. For $g\ge 2$, if we assume that the possible distributions of the random vector $\mathbf{X}=(X_1,X_2,\ldots,X_p)$ ' are continuous but not multivariate normal, the rules (2) and (3) above are, in general, no longer appropriate, and alternative procedures should be used. For example (see, e.g., chapter 6 in Seber (1984) or chapter 13 in Krzanowski (1990) for details) we could use a logistic discriminant rule or we could use the ML discriminant rule (1) replacing $\mathbf{f}_{\mathbf{j}}(\mathbf{x})$ by a nonparametric density estimator computed from $\mathbf{X}_{\mathbf{j}}$, $\mathbf{i}=1,\ldots,g$.

In this paper, we present a new method for discriminating among continuous populations which are not multivariate normal. The method is based on the multivariate generalization of the transformation of Box and Cox (1964) given by Andrews et al. (1971). In section 2, we introduce the general ideas of this new discrimination procedure. In section 3, we explicitly construct the associated discriminating rule and suggest a cross-validation method for assessing its performance. Section 4 contains an example of application and section 5 some final comments.

2. MOTIVATION

Let X be a random variable which takes values denoted by x. If X is not normal, it is often convenient to consider a transformation that might help to normalize the data. A useful family of transformations is the family of Box and Cox (1964):

$$\mathbf{x}^{(\lambda)} = \begin{cases} \frac{\mathbf{x}^{\lambda} - 1}{\lambda}, & \lambda \neq 0 \\ \log \mathbf{x}, & \lambda = 0. \end{cases}$$
 (4.1)

We assume that the following model holds

$$X^{(\lambda)} \sim N(\mu, \sigma^2),$$
 (4.2)

at least approximately. When we have a random vector $\mathbf{X}=(X_1,X_2,\ldots,X_p)'$ whose distribution is not $N_p(\mu_1,\Sigma_1)$, Andrews et al. (1971) have given the following multivariate generalization of the setup (4.1)-(4.2). We have a p-vector $\Lambda=(\lambda_1,\ldots,\lambda_p)'$ of transformation parameters, one for each dimension, such that, approximately, the model

$$X^{(\Lambda)} = (X_1^{(\lambda_1)}, X_2^{(\lambda_2)}, \dots, X_p^{(\lambda_p)}) \sim N_p(\mu, \Sigma)$$
 (5)

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holds. In both (4) and (5), it is assumed that x and the components X_j are positive. If not, x and X_j must be shifted by adding suitable constants.

Write the parameters in model (5) as $\Theta=(\Lambda,\mu,\Sigma)$. Given an observation $\mathbf{x}=(\mathbf{x}_1,\ldots,\mathbf{x}_n)$ ' from the model (5), the likelihood of \mathbf{x} is given by

$$\ell(\mathbf{x};\Theta) = (2\pi)^{-(p/2)} \left| \Sigma \right|^{-(1/2)} \exp[-(1/2)(\mathbf{x}^{(\Lambda)} - \mu)^{2} \Sigma^{-1}(\mathbf{x}^{(\Lambda)} - \mu)] J_{\Lambda}(\mathbf{x}), \tag{6}$$

where $J_{\Lambda}(\mathbf{x})$ is the jacobian of the transformation, $J_{\Lambda}(\mathbf{x}) = \int_{\mathbb{T}^{1}}^{\mathbb{P}} \mathbf{x}^{\lambda} \mathbf{j}^{-1}$. If we now face the problem of discriminating among groups G_{1}, \ldots, G_{g} such that the respective densities $f_{1}(\mathbf{x})$ are not multivariate normal but can be, in turn, modelled by (5), i.e. such that $f_{1}(\mathbf{x}) = f_{\Theta_{1}}(\mathbf{x}) = \ell(\mathbf{x}; \Theta_{1})$, where $\Theta_{1} = (\Lambda_{1}, \mu_{1}, \Sigma_{1})$, $i = 1, \ldots, g$, we can construct the following SML

discriminant rule: Assign \mathbf{x} to $\mathbf{G}_{\mathbf{i}}$ if

$$\ell(\mathbf{x}; \hat{\boldsymbol{\theta}}_{i}) = \max_{1 \le j \le g} \ell(\mathbf{x}; \hat{\boldsymbol{\theta}}_{j}), \tag{7}$$

where $\theta_i = \theta_i(x_i)$ is an efficient estimator of θ_i , i=1, ...,g. In the next section, we show how to compute explicitly the discriminant rule (7), we relate it to the rules (2) and (3), and suggest a procedure for assessing its performance.

3. CONSTRUCTION AND ASSESSMENT OF THE DISCRIMINANT RULE

3.1 Two groups

We will consider first the case when we have only two groups, G_1 and G_2 . For i=1,2, we are given a data matrix \mathfrak{X}_i of $n_i \times p$ from G_i such that for some unknown vector $\Lambda_i = (\lambda_{i1}, \ldots, \lambda_{ip})$ in \mathbb{R}^p the rows of the $n_i \times p$ transformed data matrix $\mathfrak{X}_i^{(\Lambda_i)} = (\mathfrak{X}_{ijk}^{(\Lambda_{ij})})$, j=1, ..., n_i , k=1, ...,p, are i.i.d. $N_p(\mu_i, \Sigma_i)$. We also assume independence between \mathfrak{X}_i and \mathfrak{X}_2 .

3.1.1 Parameter estimation

Standard normal likelihood theory shows (see, e.g. Gnanadesikan (1977), p. 141), that, for i=1,2, the maximum likelihood estimator (MLE) of the parameter $\Theta_i = (\Lambda_i, \mu_i, \Sigma_i)$, denoted by $\Theta_i = (\hat{\Lambda}_i, \hat{\mu}_i, \hat{\Sigma}_i)$, is given by $\hat{\mu}_i = \overline{\chi}_i^{(\hat{\Lambda}_i)}$, $\hat{\Sigma}_i = ((n_i - 1)/n_i)S_i^{(\hat{\Lambda}_i)}$, where $\overline{\chi}_i^{(\hat{\Lambda}_i)}$ and $S_i^{(\hat{\Lambda}_i)}$ are, respectively, the sample mean and the covariance matrix for $\chi_i^{(\hat{\Lambda}_i)}$. For i=1,2, $\hat{\Lambda}_i$ is obtained by maximizing the corresponding concentrated log-likelihood which is (up to an additive constant)

$$L_{\max,i}(\Lambda_i) = -(n_i/2)\log[|X_i^{(\Lambda_i)'}Q_iX_i^{(\Lambda_i)}|] + \log[J_{\Lambda_i}(X_i)], \tag{8}$$
where
$$Q_i = (I_{n_i} - 1_{n_i} 1_{n_i}' / n_i) \qquad \text{and} \qquad J_{\Lambda_i}(X_i) = \prod_{k=1}^p J_{\lambda_{ik}},$$

 $J_{\lambda_{ik}} = \prod_{j=1}^{n_1} \left| \frac{d\alpha_{ijk}^{(\lambda_{ik})}}{d\alpha_{ijk}} \right| = \left(\prod_{j=1}^{n_i} \alpha_{ijk}\right)^{\lambda_{ik}-1}, \quad k=1, \dots, p, \quad \text{are the jacobian terms. Define, for each column k of } \mathfrak{X}_i, \quad \text{the } n_i \text{ vector of normalized variables} \qquad Z_{ik}^{(\lambda_{ik})} = \left(z_{ijk}^{(\lambda_{ik})}\right), \quad k=1, \dots, p, \quad \text{of } \quad \text{jth coordinate} \\ Z_{ijk}^{(\lambda_{ik})} = \alpha_{ijk}^{(\lambda_{ik})} / J_{\lambda_{ik}}^{1/n_i}, \quad j=1, \dots, n_i. \quad \text{We have the following result.}$

LEMMA 3.1 For i=1,2, let $Z_{i}^{(\Lambda_i)} = (Z_{i1}^{(\lambda_{i1})}, ..., Z_{ip}^{(\lambda_{ip})})$ be the $n_i \times p$ matrix of normalized variables associated with \mathfrak{X}_i . We can write

$$L_{\max,i}(\Lambda_i) = -(n_i/2)\log[|Z_i^{(\Lambda_i)'}Q_iZ_i^{(\Lambda_i)}|]. \tag{9}$$

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Proof. Define, for i=1,2, the matrix $D_i = \operatorname{diag}(J_{\lambda_{1k}}^{1/n}i, \ldots, J_{\lambda_{1p}}^{1/n}i)$. (9) follows from (8), and by observing that $-(n_1/2)\log[J_{\Lambda_i}^{-2/n}i(\mathfrak{X}_i)] = -(n_1/2)\log[|D_i|^{-2}], \text{ and that } \mathbf{Z}_i^{(\Lambda_i)} = \mathfrak{X}_i^{(\Lambda_i)}D_i^{-1}.$

For computational purposes, expression (9) is more convenient than (8) in order to determine the MLE $\hat{\Lambda}_{l}$. In experience of the authors, the functions $-L_{\max,l}(\Lambda_{l})$ are typically convex, a convenient feature for solving the optimization problem by means of a canned numerical routine.

3.1.2 Expression of the rule

Given a new individual $\mathbf{x}=(\mathbf{x_1},\ldots,\mathbf{x_p})$ ' define the p vector $\mathbf{u}(\mathbf{x})=-(\log \mathbf{x_1},\ldots,\log \mathbf{x_p})$ '. Two possible different situations arise:

(i) If $\Sigma_1 \neq \Sigma_2$, our proposal is estimating $\Theta_1 = (\Lambda_1, \mu_1, \Sigma_1)$ with $(\hat{\Lambda}_1, \hat{\mu}_1, \hat{S}_1)$, where $\hat{S}_1 = S_1^{(\hat{\Lambda}_1)}$. The rule (7) can be written equivalently in the form: Assign x to G_1 if

$$\oint_{\mathbf{i}} (\mathbf{x}) + \hat{\Lambda}_{\mathbf{i}}' \mathbf{u}(\mathbf{x}) = \min_{\mathbf{1} \le \mathbf{i} \le 2} \left[\oint_{\mathbf{j}} (\mathbf{x}) + \hat{\Lambda}_{\mathbf{j}}' \mathbf{u}(\mathbf{x}) \right], \tag{10}$$

where $\hat{Q}_{1}(\mathbf{x})=(1/2)\log[|\hat{S}_{1}|]+(1/2)[(\mathbf{x}^{(\hat{\Lambda}_{1})}-\hat{\mu}_{1})'(\hat{S}_{1})^{-1}(\mathbf{x}^{(\hat{\Lambda}_{1})}-\hat{\mu}_{1})],$ i=1,2. We can interpret (10) as the estimated normal theory discriminant rule (3) applied to the transformed data matrices $\chi_{1}^{(\hat{\Lambda}_{1})}$ and $\chi_{2}^{(\hat{\Lambda}_{2})}$ and shifted by a linear term in $\mathbf{u}(\mathbf{x})$ which reflects the effect of the transformation;

(ii) If $\Sigma_1 = \Sigma_2 = \Sigma$, it is reasonable to estimate Σ by the pooled estimator $\hat{S}_p = [(n_1 - 1)\hat{S}_1 + (n_2 - 1)\hat{S}_2]/(N-2)$, where $N = n_1 + n_2$. The rule (7) becomes now: Assign \mathbf{x} to G_1 if

$$\widetilde{Q}_{j}(\mathbf{x}) + \widehat{\Lambda}_{j}' \mathbf{u}(\mathbf{x}) = \min_{1 \le j \le 2} [\widetilde{Q}_{j}(\mathbf{x}) + \widehat{\Lambda}_{j}' \mathbf{u}(\mathbf{x})], \tag{11}$$

where
$$\tilde{Q}_{i}(\mathbf{x})=(1/2)[(\mathbf{x}^{(\hat{\Lambda}_{i})}-\hat{\mu}_{i}),\hat{S}_{p}^{-1}(\mathbf{x}^{(\hat{\Lambda}_{i})}-\hat{\mu}_{i})], i=1,2.$$

A criterion for deciding whether (10) or (11) should be used is suggested in the analysis of the example in section 4 below.

3.1.3 Cross-validation assessment

To evaluate the performance of the rules (10) and (11), we need to compute estimates of the error rates p_{ij} =P[Assign x to $G_i | G_j |$, i.e., the probabilities of allocating an individual to G_i when, in fact, it comes from G_j . A well established method for estimating the p_{ij} is the method of cross-validation proposed by Lachenbruch and Mickey (1968). The general idea in cross-validation is determining the allocation rule using the training data deleting one observation, and then using the rule to classify the observation left out. The estimates are $\hat{P}_{ij} = a_{ij}/n_i$, where a_{ij} is the number of observations from G_i which have been misclassified.

For the rules (10) and (11), we need to compute $\hat{\Lambda}_{i(j)}$, the MLE of Λ_i deleting the jth row of the data matrix \mathfrak{X}_i , i=1,2, j=1, ..., n_i . This is computationally expensive and, therefore, an approximation is in order.

By (9), the MLE of Λ_i is obtained by minimizing $|M_i(\Lambda_i)|$, where $M_i(\Lambda_i) = Z_i^{(\Lambda_i)'}Q_iZ_i^{(\Lambda_i)}$. $\hat{\Lambda}_{i(j)}$ minimizes $|M_{i(j)}(\Lambda_i)|$. Let $\Lambda_{i,0}$ an initial guess for $\hat{\Lambda}_{i(j)}$ (typically $\hat{\Lambda}_i$) and let $\nabla_{i(j)}(\Lambda_i) = \partial |M_{i(j)}(\Lambda_i)| / \partial \Lambda_i$ and $H_{i(j)}(\Lambda_i) = \partial^2 |M_{i(j)}(\Lambda_i)| / \partial \Lambda_i \partial \Lambda_i'$, be, respectively, the gradient vector and the hessian matrix of $|M_{i(j)}(\Lambda_i)|$. In the first-order Taylor expansion $\nabla_{i(j)}(\Lambda_i) \cong \nabla_{i(j)}(\Lambda_{i,0}) + H_{i(j)}(\Lambda_{i,0}) (\Lambda_i - \Lambda_{i,0})$, we can use the fact that $\nabla_{i(j)}(\hat{\Lambda}_{i(j)}) = 0$ to obtain a one-step approximation for $\hat{\Lambda}_{i(j)}$,

$$\hat{\Lambda}_{1(j)}^{1} = \Lambda_{1,0} - (H_{1(j)}(\Lambda_{1,0}))^{-1} \nabla_{1(j)}(\Lambda_{1,0}). \tag{12}$$

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Equation (12) is a multivariate extension of equation (15) in Tsai and Wu (1990).

To compute $\nabla_{i(j)}(\hat{\Lambda}_i)$ and $H_{i(j)}(\hat{\Lambda}_i)$ in (12), we define, for i=1,2; $j=1,\ldots,n_i$ and k=1, ...,p: (i) $a_{ijk}=x_{ijk}^{1/n}i(\prod_{m\neq j}x_{imk})^{-1/n}i^{(n_i-1)},$ $q_{ijk}(\lambda_{ik})=(a_{ijk})^{\lambda_{ik}-1};$ (ii) $Z_{ijk}^{\bullet(\lambda_{ik})}=q_{ijk}(\lambda_{ik})Z_{ik}^{(\lambda_{ik})};$ and (iii) $A_{ij}=Q_iB_{ij}Q_i$, $B_{ij}=I_n-[1-(1/n_i)]^{-1}e_{ij}e_{ij}$, where e_{ij} is the jth canonical vector of \mathbb{R}^n . We have the following result.

THEOREM 3.1 For i=1,2 and j=1, ..., n_i , define the $n_i \times p$ matrix $Z_{ij}^{*(\Lambda_i)} = (Z_{ij1}^{*(\lambda_{ij1})}, \ldots, Z_{ijp}^{*(\lambda_{ip})})$. We can write:

$$\mathbf{M}_{\mathbf{i}(\mathbf{j})}(\Lambda_{\mathbf{i}}) = \mathbf{Z}_{\mathbf{i}\,\mathbf{j}}^{\bullet(\Lambda_{\mathbf{i}})} \mathbf{A}_{\mathbf{i}\,\mathbf{j}} \mathbf{Z}_{\mathbf{i}\,\mathbf{j}}^{\bullet(\Lambda_{\mathbf{i}})}. \tag{13}$$

Proof. We will use the following formula (see Atkinson (1986, p.31)): If a,b are two vectors in \mathbb{R}^n and $P=(p_{uv})$ is an nxn orthogonal projection matrix, we have, for $l=1, \ldots, n$,

$$a'(I_{n}-P)b=a'_{(1)}(I_{n-1}-P_{(1)})b_{(1)}+(1-p_{11})^{-1}a_{1}^{*}b_{1}^{*}, \qquad (14)$$

The subscript (1) means that the corresponding quantity has been computed deleting the 1th row or coordinate. a_1^* and b_1^* are the 1th coordinates of the residual vectors $a = (I_n - P)a$ and $b = (I_n - P)b$. For

i=1,2, j=1, ..., n_i and r,s=1, ...,p, the (r,s) element in the pxp matrix $M_{I(j)}(\Lambda_i) = (m_{Irs(j)}(\lambda_{Ir},\lambda_{Is}))$ is $m_{Irs(j)} = Z_{Ir}^{(\lambda_{Ir})'}Q_{I(j)}Z_{Is}^{(\lambda_{Is})}$. By using an analog of expression (13) in Tsai and Wu (1990) and formula (14) above, we get $m_{Irs(j)} = Z_{Ijr}^{*(\lambda_{Ir})}A_{Ij}Z_{Ijs}^{*(\lambda_{Is})}$.

It is important to observe that, by (13), $\hat{\Lambda}_{l(j)}$ minimizes a determinant which is structurally similar as the determinant which minimizes $\hat{\Lambda}_{l}$. Expressions for the partial derivatives of $|M_{l(j)}(\Lambda_{l})|$ can be seen in the appendix.

3.2 More than two groups

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When g>2, the construction and assessment of the discriminant rule (7) can be done following straightforward extensions of the ideas in sections 3.1.1, 3.1.2 and 3.1.3.

4. A PRACTICAL EXAMPLE

In DELPHI, one of the detectors of LEP, a particle collider at the European Center for Nuclear Research (CERN, Geneva, Switzerland), several variables are observed by measuring the events produced in the particle collisions. The problem is classifying, or "tagging", on the basis of the observed measurements, the event containing a quark into one of three different groups:

G: Up and Down (u+d) or Strange (s);

G₂: Charm (c);

 G_3 : Beauty (b).

Among the set of variables which are observed, we will consider only a special subset of 9 continuous variables (X_1, \ldots, X_9) , called Microvertex variables. We are given simulated data matrices \mathfrak{X}_1 from the

groups G_1 , such that \mathfrak{X}_1 is $n_1 \times 9$, with $n_1 = 182$, $n_2 = 52$, and $n_3 = 66$. We illustrate, with reference to this example, the discriminant techniques presented in section 3. The problem of classification has been also treated by de la Vaissière and Palma-Lopes (1989), two physicists at CERN.

The use of the rule in section 3 arises because of the strong nonnormality of the data. For example, figure 1 below shows the skewed histograms for the variable X_4 for the groups G_1 , G_2 and G_3 . We decided then, applying the multivariate transformation given in (5). Since, for other variables, negative data appear, we consider, for i=1,2,3, j=1, k=1, ...,9, $x_{lik} \longrightarrow x_{lik} + a_{k}$ translation the a = -x +.5 and x is the global minimum of the kth variable over the three data matrices \mathfrak{X}_1 , \mathfrak{X}_2 , \mathfrak{X}_3 , k=1, ...,p. We could have considered a shifted version of the transformation (4). Unfortunately, shifting the transformation (4) causes the maximum likelihood estimation of the transformation parameters to become a non-regular problem and not a well established procedure exists yet. However, see a recent work of Atkinson et al. (1991).

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Figure 1

By maximizing the respective functions (9) for each class, we found the set of estimated transformation parameters displayed in table 1 below.

Table 1

In order to use (10) or (11) we need now to decide whether or not the dispersion matrices of the transformed data are the same. A reasonable

procedure is to test for homoscedasticity treating the transformed data matrices $\mathfrak{A}_1^{(\hat{\Lambda})}$, $\mathfrak{A}_2^{(\hat{\Lambda})}$ and $\mathfrak{A}_3^{(\hat{\Lambda})}$ as data matrices from a multivariate normal distribution. The standard likelihood ratio test, as described for example in Mardia et al. (1979), p. 140, revealed strong evidence again homoscedasticity and, therefore, we decided to use the rule of section 3 in the form (10). It is important to remark that these procedure for choosing between (10) and (11) is not entirely rigorous because we are ignoring the effect of the sampling variability of the $\hat{\Lambda}_1$ in the null distribution of the likelihood ratio test statistic for homoscedasticity, i.e. we are treating the $\hat{\Lambda}_1$ as known constants. This procedure should be therefore taken as merely an indicative guideline for choosing between (10) and (11).

Finally, table 2 shows the estimated error rates by cross-validation. The rule classifies correctly slightly less than 70% of the events containing a beauty quark. This is remarkable because, from a physical point of view, these are the events where correct classification is more important. The relatively poor performance of the rule regarding events containing a quark of type c is not unexpected because G_2 is a transition class between G_1 and G_3 and, therefore, not very well differentiated from G_1 and G_2 .

Table 2

5. CONCLUSIONS

This paper presents a new discriminant rule for the case when the groups do not have a multivariate normal distribution but can be

modelled approximately as normal after a suitable transformation. From a practical point of view, the rule presented is particularly useful in the presence of long tailed asymmetric marginal distributions which are susceptible of treatment under the framework (4.1)-(4.2). The implementation of the rule should be accompanied with a diagnostic check of normality and heteroscedasticity of the transformed data.

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APPENDIX: ELEMENTS OF
$$\nabla_{1(j)}(\hat{\Lambda}_{1})$$
 AND $H_{1(j)}(\hat{\Lambda}_{1})$.

For a general nxp matrix $H=(h_1,\ldots,h_p)$, we will write $H_j(u)$ for the nxp matrix obtained from H by replacing its jth column by the nxl vector u. We will also write $H_{jl}(u,v)$ for the nxp matrix obtained from H by replacing its jth and lth column by, respectively, u and v.

THEOREM A.1 Write, for i=1,2, $\hat{\Lambda}_i = (\hat{\lambda}_{i1}, \ldots, \hat{\lambda}_{ip})$ ' and define, for i=1,2, j=1, ...,p, and k=1, ...,p, the $n_i x l$ vectors $Z_{ijk}^{\bullet(\lambda_{ik})} = q_{ijk}(\lambda_{ik}) Z_{ik}^{(\lambda_{ik})}$ (see section 3.1.3), $W_{ijk}^{\bullet(\lambda_{ik})} = \partial Z_{ijk}^{\bullet(\lambda_{ik})} / \partial \lambda_{ik}$ and $U_{ijk}^{\bullet(\lambda_{ik})} = \partial^2 Z_{ijk}^{\bullet(\lambda_{ik})} / \partial \lambda_{ik}^2$. We will write Z_{ijk}^{\bullet} , W_{ijk}^{\bullet} and U_{ijk}^{\bullet} for the corresponding vectors evaluated at $\hat{\lambda}_{ik}$. If $Z_{ij}^{\bullet} = Z_{ij}^{\bullet(\hat{\Lambda}_i)}$ (see section 3.1.3), we have, for r,s=1, ...,p,

(a)
$$\partial \left| M_{i(j)}(\Lambda_i) \right| / \partial \lambda_{ir} \left| \Lambda_i = \hat{\Lambda}_i = 2 \left| \mathbf{Z}_{ij,r}^{\bullet}(\mathbf{W}_{ijr}^{\bullet}) \mathbf{A}_{ij} \mathbf{Z}_{ij}^{\bullet} \right|;$$

(b)
$$\partial^{2} |M_{i(j)}(\Lambda_{i})| / \partial^{2} \lambda_{ir} |_{\Lambda_{i} = \hat{\Lambda}_{i}} = 2[|Z_{ij,r}^{\bullet}(W_{ijr}^{\bullet})A_{ij}Z_{ij,r}^{\bullet}(W_{ijr}^{\bullet})| + |Z_{ij}^{\bullet}A_{ij}Z_{ij,r}^{\bullet}(U_{ijr}^{\bullet})|];$$

(c)
$$\partial^{2} |M_{1(j)}(\Lambda_{i})| / \partial \lambda_{ir} \partial \lambda_{is} |_{\Lambda_{i} = \hat{\Lambda}_{i}} = 2 [|Z_{ij,rs}^{*}(W_{ijr}^{*}, W_{ijs}^{*}) A_{ij} Z_{ij}^{*}| + |Z_{ij,r}^{*}(W_{ijr}^{*}) A_{ij} Z_{ij,s}^{*}(W_{ijs}^{*})|]. (r \neq s).$$

To obtain these expressions we need the following lemma.

LEMMA A.1 Let C, D, and B be three nxp matrices and let E be a symmetric matrix of nxn. If e_i represents the ith canonical vector of \mathbb{R}^n , we have, for i=1, ...,n, j=1, ...,p, the following differentiation formulae: a) $\partial |C'ED|/\partial c_{ij} = |C'_j(e_i)ED|$; b) $\partial |C'ED|/\partial d_{ij} = |C'ED_j(e_i)|$; and c) $\partial |B'EB|/\partial b_{ij} = 2|B'_j(e_i)EB|$.

Proof. See lemma 3.1 in Velilla (1992).

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Proof of theorem A.1. By the chain rule and part a) of lemma A.1, $\frac{\partial \left| M_{i(j)}(\Lambda_i) \right|}{\partial \lambda_{ir}} \left| \frac{\partial \left| \Lambda_i \right|}{\Lambda_i} \right| = \sum_{i=1}^n \sum_{m=1}^p \left| Z_{ij,m}^*(e_{ii}) A_{ij} Z_{ij}^* \right| \left(\frac{\partial Z_{ij,im}^*(\lambda_i)}{\partial \lambda_{ir}} \right| \frac{\partial A_{ir}}{\partial \lambda_{ir}} \right| = 2 \sum_{i=1}^n \left| Z_{ij,r}^*(W_{ijr}^*) A_{ij} Z_{ij}^* \right|.$ Expressions b) and c) are obtained in a similar way by using, respectively, parts a) and b) of lemma A.1.

Finally, if from $Z_{ik}^{(\lambda)}_{ik}$ (see section 3.1.1), we define $W_{ik}^{(\lambda)}_{ik}=\partial Z_{ik}^{(\lambda)}_{ik}/\partial \lambda_{ik}$, $U_{ik}^{(\lambda)}_{ik}=\partial^2 Z_{ik}^{(\lambda)}_{ik}/\partial \lambda_{ik}^2$ and Z_{ik} , W_{ik} , U_{ik} stand for the corresponding functions evaluated at $\hat{\lambda}_{ik}$, it is easily seen that $W_{ijk}^*=q_{ijk}(\hat{\lambda}_{ik})[\log(a_{ijk})Z_{ik}+W_{ik}]$ and $U_{ijk}^*=q_{ijk}(\hat{\lambda}_{ik})[(\log(a_{ijk}))^2Z_{ik}+2\log(a_{ijk})W_{ik}+U_{ik}]$. As a conclusion, for the cross-validation assessment of the rules (10) and (11), we need the two MLE's $\hat{\lambda}_{i}$ and $\hat{\lambda}_{i}$, the array of constants (a_{ijk}) and the n_{ik} vectors Z_{ik} , W_{ik} and U_{ik} . An explicit expression for the functions $Z_{ik}^{(\lambda)}_{ik}$, $W_{ik}^{(\lambda)}_{ik}$ and $U_{ik}^{(\lambda)}_{ik}$, can be found in Atkinson and Lawrance (1989).

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CAPTIONS FOR TABLES AND FIGURES

Table 1. Estimated transformation parameters for groups G_1 , G_2 and G_3 .

Table 2. Estimated error rates by cross-validation.

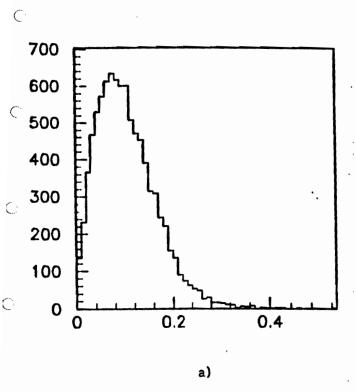
Figure 1. Histograms for X_4 for groups G_1 , G_2 and G_3 . a) G_1 ; b) G_2 ; c) G_3 .

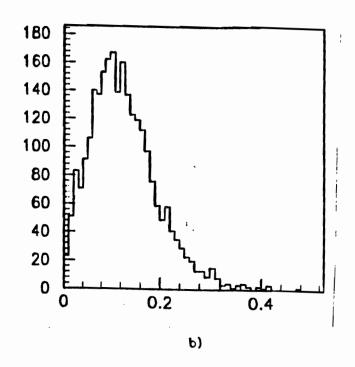
	1	2	3	4	5	6	7	8	9
۾	0.985	1.437	0.108	0.485	-2.534	-1.347	1.185	1.025	1.245
۸,	0.295	-0.313	0.550	0.558	-1.736	-1.043	1.062	0.853	0.826
۲,	-0.622	0.405	0.930	0.582	-0.922	-0.377	0.303	0.695	0.976

Table 1

	G ₁	$G_{_{2}}$	
G ₁	0.621	0.236	0.143
G ₂	0.481	0.231	0.288
G ₃	0.182	0.136	0.682

Table 2





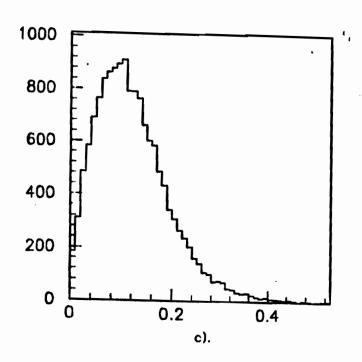


Figure 1