# A Tutorial on Multivariate <br> Statistical Analysis 

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## ELEMENTARY STATISTICS

Collection of (real-valued) data from a sequence of experiments

$$
X_{1}, X_{2}, \ldots, X_{n}
$$

Might make assumption underlying law is $N\left(\mu, \sigma^{2}\right)$ with unknown mean $\mu$ and variance $\sigma^{2}$. Want to estimate $\mu$ and $\sigma^{2}$ from the data. Sample Mean \& Sample Variance:

$$
\bar{X}=\frac{1}{n} \sum_{j} X_{j}, \quad S=\frac{1}{n-1} \sum_{j}\left(X_{j}-\bar{X}\right)^{2}
$$

Estimators are "unbiased"

$$
\mathbb{E}(\bar{X})=\mu, \quad \mathbb{E}(S)=\sigma^{2}
$$

Theorem: If $X_{1}, X_{2}, \ldots$ are independent $N\left(\mu, \sigma^{2}\right)$ variables then $\bar{X}$ and $S$ are independent. We have that $\bar{X}$ is $N\left(\mu, \sigma^{2} / n\right)$ and $(n-1) S / \sigma^{2}$ is $\chi^{2}(n-1)$.

Recall $\chi^{2}(d)$ denotes the chi-squared distribution with $d$ degrees of freedom. Its density is

$$
f_{\chi^{2}}(x)=\frac{1}{2^{d / 2} \Gamma(d / 2)} x^{d / 2-1} \mathrm{e}^{-x / 2}, \quad x \geq 0
$$

where

$$
\Gamma(z)=\int_{0}^{\infty} t^{z-1} \mathrm{e}^{-t} d t, \Re(z)>0
$$

## MULTIVARIATE GENERALIZATIONS

From the classic textbook of Anderson[1]:
Multivariate statistical analysis is concerned with data that consists of sets of measurements on a number of individuals or objects. The sample data may be heights and weights of some individuals drawn randomly from a population of school children in a given city, or the statistical treatment may be made on a collection of measurements, such as lengths and widths of petals and lengths and widths of sepals of iris plants taken from two species, or one may study the scores on batteries of mental tests administered to a number of students.

$$
\begin{aligned}
& p=\# \text { of sets of measurements on a given individual, } \\
& n=\# \text { of observations }=\text { sample size }
\end{aligned}
$$

## Remarks:

- In above examples, one can assume that $p \ll n$ since typically many measurements will be taken.
- Today it is common for $p \gg 1$, so $n / p$ is no longer necessarily large.

Vehicle Sound Signature Recognition: Vehicle noise is a stochastic signal. The power spectrum is discretized to a vector of length $p=1200$ with $n \approx 1200$ samples from the same kind of vehicle.

Astrophysics: Sloan Digital Sky Survey typically has many observations (say of quasar spectrum) with the spectra of each quasar binned resulting in a large $p$.

Financial data: S\&P 500 stocks observed over monthly intervals for twenty years.

## GAUSSIAN DATA MATRICES

The data are now $n$ independent column vectors of length $p$

$$
\vec{x}_{1}, \vec{x}_{2}, \ldots, \vec{x}_{n}
$$

from which we construct the $n \times p$ data matrix

$$
X=\left(\begin{array}{ccc}
\longleftarrow & \vec{x}_{1}^{T} & \longrightarrow \\
\longleftarrow & \vec{x}_{2}^{T} & \longrightarrow \\
& \vdots & \\
\longleftarrow & \vec{x}_{n}^{T} & \longrightarrow
\end{array}\right)
$$

The Gaussian assumption is that

$$
\vec{x}_{j} \sim N_{p}(\mu, \Sigma)
$$

Many applications assume the mean has been already substracted out of the data, i.e. $\mu=0$.

## Multivariate Gaussian Distribution

If $x$ and $y$ are vectors, the matrix $x \otimes y$ is defined by

$$
(x \otimes y)_{j k}=x_{j} y_{k}
$$

If $\mu=\mathbb{E}(x)$ is the mean of the random vector $x$, then the covariance matrix of $x$ is the $p \times p$ matrix

$$
\Sigma=\mathbb{E}[(x-\mu) \otimes(x-\mu)]]
$$

$\Sigma$ is a symmetric, non-negative definite matrix. If $\Sigma>0$ (positive definite) and $X \sim N_{p}(\mu, \Sigma)$, then the density function of $X$ is

$$
f_{X}(x)=(2 \pi)^{-p / 2}(\operatorname{det} \Sigma)^{-1 / 2} \exp \left[-\frac{1}{2}\left(x-\mu, \Sigma^{-1}(x-\mu)\right],>x \in \mathbb{R}^{p}\right.
$$

Sample mean:

$$
\bar{x}=\frac{1}{n} \sum_{j} \vec{x}_{j}, \quad \mathbb{E}(\bar{x})=\mu
$$

Sample covariance matrix:

$$
S=\frac{1}{n-1} \sum_{j=1}^{n}\left(\vec{x}_{j}-\bar{x}\right) \otimes\left(\vec{x}_{j}-\bar{x}\right)
$$

For $\mu=0$ the sample covariance matrix can be written simply as

$$
\frac{1}{n-1} X^{T} X
$$

Some Notation: If $X$ is a $n \times p$ data matrix formed from the $n$ independent column vectors $x_{j}, \operatorname{cov}\left(x_{j}\right)=\Sigma$, we can form one column $\operatorname{vector} \operatorname{vec}(X)$ of length $p n$

$$
\operatorname{vec}(X)=\left(\begin{array}{l}
x_{1} \\
x_{2} \\
\vdots \\
x_{n}
\end{array}\right)
$$

The covariance of $\operatorname{vec}(X)$ is the $n p \times n p$ matrix

$$
I_{n} \otimes \Sigma=\left(\begin{array}{ccccc}
\Sigma & 0 & 0 & \cdots & 0 \\
0 & \Sigma & 0 & \cdots & 0 \\
\vdots & \vdots & \vdots & \cdots & \vdots \\
0 & 0 & 0 & \cdots & \Sigma
\end{array}\right)
$$

In this case we say the data matrix $X$ constructed from $n$ independent $x_{j} \sim N_{p}(\mu, \Sigma)$ has distribution

$$
N_{p}\left(M, I_{n} \otimes \Sigma\right)
$$

where $M=\mathbb{E}(X)=\mathbf{1} \otimes \mu, \mathbf{1}$ is the column vector of all 1 's.

## WISHART DISTRIBUTION

Definition: If $A=X^{T} X$ where the $n \times p$ matrix $X$ is $N_{p}\left(0, I_{n} \otimes \Sigma\right), \Sigma>0$, then $A$ is said to have Wishart distribution with $n$ degrees of freedom and covariance matrix $\Sigma$. We will say $A$ is $W_{p}(n, \Sigma)$.

## Remarks:

- The Wishart distribution is the multivariate generalization of the chi-squared distribution.
- $A \sim W_{p}(n, \Sigma)$ is positive definite with probability one if and only if $n \geq p$.
- The sample covariance matrix,

$$
S=\frac{1}{n-1} A
$$

is $W_{p}\left(n-1, \frac{1}{n-1} \Sigma\right)$.

## WISHART DENSITY FUNCTION, $n \geq p$

Let $\mathcal{S}_{p}$ denote the space of $p \times p$ positive definite (symmetric) matrices. If $A=\left(a_{j k}\right) \in \mathcal{S}_{p}$, let

$$
(d A)=\text { volume element of } A=\bigwedge_{j \leq k} d a_{j k}
$$

The multivariate gamma function is

$$
\Gamma_{p}(a)=\int_{\mathcal{S}_{p}} \mathrm{e}^{-\operatorname{tr}(A)}(\operatorname{det} A)^{a-(p+1) / 2}(d A), \Re(a)>(p-1) / 2 .
$$

Theorem: If $A$ is $W_{p}(n, \Sigma)$ with $n \geq p$, then the density function of $A$ is

$$
\frac{1}{2^{n p} \Gamma_{p}(n / 2)(\operatorname{det} \Sigma)^{n / 2}} \mathrm{e}^{-\frac{1}{2} \operatorname{tr}\left(\Sigma^{-1} A\right)}(\operatorname{det} A)^{(n-p-1) / 2}
$$

## Sketch of Proof:

- The density function for $X$ is the multivariate Gaussian (including volume element ( $d X$ ))

$$
(2 \pi)^{-n p / 2}(\operatorname{det} \Sigma)^{-n / 2} \mathrm{e}^{-\frac{1}{2} \operatorname{tr}\left(\Sigma^{-1} X^{T} X\right)}(d X)
$$

- Recall the $Q R$ factorization [8]: Let $X$ denote an $n \times p$ matrix with $n \geq p$ with full column rank. Then there exists an unique $n \times p$ matrix $Q, Q^{T} Q=I_{p}$, and an unique $n \times p$ upper triangular matrix $R$ with positive diagonal elements so that $X=Q R$. Note $A=X^{T} X=R^{T} R$.
- A Jacobian calculation $[1,13]$ : If $A=X^{T} X$, then

$$
(d X)=2^{-p}(\operatorname{det} A)^{(n-p-1) / 2}(d A)\left(Q^{T} d Q\right)
$$

where

$$
\left(Q^{T} d Q\right)=\bigwedge_{j=1}^{p} \bigwedge_{k=j+1}^{n} q_{k}^{T} d q_{j}
$$

and $Q=\left(q_{1}, \ldots, q_{p}\right)$ is the column representation of $Q$.

- Thus the joint distribution of $A$ and $Q$ is

$$
\begin{aligned}
& (2 \pi)^{-n p / 2}(\operatorname{det} \Sigma)^{-n / 2} \mathrm{e}^{-\frac{1}{2} \operatorname{tr}\left(\Sigma^{-1} A\right)} \times \\
& 2^{-p}(\operatorname{det} A)^{(n-p-1) / 2}(d A)\left(Q^{T} d Q\right)
\end{aligned}
$$

- Now integrate over all $Q$. Use fact that

$$
\int_{\mathcal{V}_{n, p}}\left(Q^{T} d Q\right)=\frac{2^{p} \pi^{n p / 2}}{\Gamma_{p}(n / 2)}
$$

and $\mathcal{V}_{n, p}$ is the set of real $n \times p$ matrices $Q$ satisfying $Q^{T} Q=I_{p}$. (When $n=p$ this is the orthogonal group.)

Remarks regarding the Wishart density function

- Case $p=2$ obtain by R. A. Fisher in 1915.
- General $p$ by J. Wishart in 1928 by geometrical arguments.
- Proof outlined above came later. (See [1, 13] for complete proof.)
- When $Q$ is a $p \times p$ orthogonal matrix

$$
\frac{\Gamma_{p}(p / 2)}{2^{p} \pi^{p^{2} / 2}}\left(Q^{T} d Q\right)
$$

is normalized Haar measure for the orthogonal group $\mathcal{O}(p)$. We denote this Haar measure by $(d Q)$.

- Siegel proved (see, e.g. [13])

$$
\Gamma_{p}(a)=\pi^{p(p-1) / 4} \prod_{j=1}^{p} \Gamma\left(a-\frac{1}{2}(j-1)\right)
$$

## EIGENVALUES OF A WISHART MATRIX

Theorem: If $A$ is $W_{p}(n, \Sigma)$ with $n \geq p$ the joint density function for the eigenvalues $\ell_{1}, \ldots, \ell_{p}$ of $A$ is

$$
\begin{array}{r}
\frac{\pi^{p^{2} / 2} 2^{-n p / 2}(\operatorname{det} \Sigma)^{-n / 2}}{\Gamma_{p}(p / 2) \Gamma_{p}(n / 2)} \prod_{j=1}^{p} \ell_{j}^{(n-p-1) / 2} \prod_{j<k}\left|\ell_{j}-\ell_{k}\right| \times \\
\int_{\mathcal{O}(p)} \mathrm{e}^{-\frac{1}{2} \operatorname{tr}\left(\Sigma^{-1} Q L Q^{T}\right)}(d Q), \quad\left(\lambda_{1}>\cdots>\lambda_{p}\right)
\end{array}
$$

where $L=\operatorname{diag}\left(\ell_{1}, \ldots, \ell_{p}\right)$ and $(d Q)$ is normalized Haar measure. Note that $\Delta(\ell):=\prod_{j<k}\left(\ell_{j}-\ell_{k}\right)$ is the Vandermonde.
Corollary: If $A$ is $W_{p}\left(n, I_{p}\right)$, then the integral over the orthogonal group in the previous theorem is

$$
\mathrm{e}^{-\frac{1}{2} \sum_{j} \ell_{j}}
$$

Proof: Recall that the Wishart density function (times the volume element) is

$$
\frac{1}{2^{n p} \Gamma_{p}(n / 2)(\operatorname{det} \Sigma)^{n / 2}} \mathrm{e}^{-\frac{1}{2} \operatorname{tr}\left(\Sigma^{-1} A\right)}(\operatorname{det} A)^{(n-p-1) / 2}(d A)
$$

The idea is to diagonalize $A$ by an orthogonal transformation and then integrate over the orthogonal group thereby giving the density function for the eigenvalues of $A$.
Let $\ell_{1}>\cdots>\ell_{p}$ be the ordered eigenvalues of $A$.

$$
A=Q L Q^{T}, \quad L=\operatorname{diag}\left(\ell_{1}, \ldots, \ell_{p}\right), \quad Q \in \mathcal{O}(p)
$$

The $j^{\text {th }}$ column of $Q$ is a normalized eigenvector of $A$. The transformation is not $1-1$ since $Q=\left[ \pm q_{1}, \ldots, \pm q_{p}\right]$ works for each fixed $A$. The transformation is made $1-1$ by requiring that the $1^{s t}$ element of each $q_{j}$ is nonnegative. This restricts $Q$ (as $A$ varies) to a $2^{-p}$ part of $\mathcal{O}(p)$. We compensate for this at the end.

We need an expression for the volume element $(d A)$ in terms of $Q$ and $L$. First we compute the differential of $A$

$$
\begin{aligned}
d A & =d Q L Q^{T}+Q d L Q^{T}+Q L d Q^{T} \\
Q^{T} d A Q & =Q^{T} d Q L+d L+L d Q^{T} Q \\
& =-d Q^{t} Q d L+L d Q^{T} Q+d L \\
& =\left[L, d Q^{T} Q\right]+d L
\end{aligned}
$$

(We used $Q^{T} Q=I$ implies $Q^{T} d Q=-d Q^{T} Q$.)
We now use the following fact (see, e.g., page 58 in [13]): If $X=B Y B^{T}$ where $X$ and $Y$ are $p \times p$ symmetric matrices, $B$ is a nonsingular $p \times p$ matrix, then $(d X)=(\operatorname{det} B)^{p+1}(d Y)$. In our case $Q$ is orthogonal so the volume element $(d A)$ equals the volume element $\left(Q^{T} d A Q\right)$. The volume element is the exterior product of the diagonal elements of $Q^{T} d A Q$ times the exterior product of the elements above the diagonal.

Since $L$ is diagonal, the commutator

$$
\left[L, d Q^{T} Q\right]
$$

has zero diagonal elements. Thus the exterior product of the diagonal elements of $Q^{T} d A Q$ is $\bigwedge_{j} d \ell_{j}$.
The exterior product of the elements coming from the commutator is

$$
\prod_{j<k}\left(\ell_{j}-\ell_{k}\right) \bigwedge_{j<k} q_{k}^{T} d q_{j}
$$

and so

$$
\begin{aligned}
(d A) & =\bigwedge_{j<k} q_{k}^{T} d q_{j} \Delta(\ell) \bigwedge_{j} d \ell_{j} \\
& =\frac{2^{p} \pi^{p^{2} / 2}}{\Gamma_{p}(p / 2)}(d Q) \Delta(\ell) \bigwedge_{j} d \ell_{j}
\end{aligned}
$$

The theorem now follows once integrate over all of $\mathcal{O}(p)$ and divide the result by $2^{p}$.

- One is interested in limit laws as $n, p \rightarrow \infty$. For $\Sigma=I_{p}$, Johnstone [11] proved, using RMT methods, for centering and scaling constants

$$
\begin{aligned}
& \mu_{n p}=(\sqrt{n-1}+\sqrt{p})^{2} \\
& \sigma_{n p}=(\sqrt{n-1}+\sqrt{p})\left(\frac{1}{\sqrt{n-1}}+\frac{1}{\sqrt{p}}\right)^{1 / 3}
\end{aligned}
$$

that

$$
\frac{\ell_{1}-\mu_{n p}}{\sigma_{n p}}
$$

converges in distribution as $n, p \rightarrow \infty, n / p \rightarrow \gamma<\infty$, to the GOE largest eigenvalue distribution [15].

- El Karoui [6] has extended the result to $\gamma \leq \infty$. The case $p \gg n$ appears, for example, in microarray data.
- Soshnikov [14] has lifted Gaussian assumption under the additional restriction $n-p=\mathrm{O}\left(p^{1 / 3}\right)$.
- For $\Sigma \neq I_{p}$, the difficulty in establishing limit theorms comes from the integral

$$
\int_{\mathcal{O}(p)} \mathrm{e}^{-\frac{1}{2} \operatorname{tr}\left(\Sigma^{-1} Q \Lambda Q^{T}\right)}(d Q)
$$

Using zonal polynomials infinite series expansions have been derived for this integral, but these expansions are difficult to analyze. See Muirhead [13].

- For complex Gaussian data matrices $X$ similar density formulas are known for the eigenvalues of $X^{*} X$. Limit theorems for $\Sigma \neq I_{p}$ are known since the analogous group integral, now over the unitary group, is known explicitly-the Harish Chandra-Itzykson-Zuber (HCIZ) integral (see, e.g. [17]). See the work of Baik, Ben Arous and Péché [2, 3] and El Karoui [7].


## PRINCIPAL COMPONENT ANALYSIS (PCA),

## H. Hotelling, 1933

Population Principal Components: Let $x$ be a $p \times 1$ random vector with $\mathbb{E}(x)=\mu$ and $\operatorname{cov}(x)=\Sigma>0$. Let $\lambda_{1} \geq \lambda_{2} \geq \cdots \geq \lambda_{p}$ denote the eigenvalues of $\Sigma$ and $H$ an orthogonal matrix diagonalizing $\Sigma: H^{T} \Sigma H=\Lambda=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{p}\right)$. We write $H$ in column vector form

$$
H=\left[h_{1}, \ldots, h_{p}\right]
$$

so that $h_{j}$ is the $p \times 1$ eigenvector of $\Sigma$ corresponding to eigenvalue $\lambda_{j}$. Define the $p \times 1$ vector

$$
u=H^{T} x=\left(u_{1}, \ldots, u_{p}\right)^{T}
$$

then

$$
\begin{aligned}
\operatorname{cov}(u) & =\mathbb{E}\left(\left(H^{T} x-H^{T} \mu\right) \otimes\left(H^{T} x-H^{T} \mu\right)\right) \\
& =H^{T} \mathbb{E}((x-\mu) \otimes(x-\mu)) H \\
& =H^{T} \Sigma H=\Lambda \Rightarrow u_{j} \text { uncorrelated. }
\end{aligned}
$$

Definition: $u_{j}$ is called the $j^{t h}$ principal component of $x$. Note $\operatorname{var}\left(u_{j}\right)=\lambda_{j}$.
Statistical interpretations: The claim is that $u_{1}$ is that linear combination of components of $x$ that has maximum variance.
Proof: For simplicity of notation, set $\mu=0$. Let $b$ denote any $p \times 1$ vector, $b^{T} b=1$, and form $b^{T} x$.
$\operatorname{var}\left(b^{T} x\right)=\mathbb{E}\left(b^{T} x \cdot b^{T} x\right)=\mathbb{E}\left(b^{T} x \cdot\left(b^{T} x\right)^{T}\right)=b^{T} \mathbb{E}\left(x x^{T}\right) b=b^{T} \Sigma b$.
We want to maximize the right hand side subject to the constraint $b^{T} b=1$. By the method of Lagrange multipliers we maximize

$$
b^{T} \Sigma b-\lambda\left(b^{T} b-1\right)
$$

Since $\Sigma$ is symmetric the vector of partial derivatives is

$$
2 \Sigma b-2 \lambda b
$$

Thus $b$ must be an eigenvector with eigenvalue $\lambda$.

The largest variance corresponds to choosing the largest eigenvalue.
The general result is that $u_{r}$ has maximum variance of all normalized combinations uncorrelated with $u_{1}, \ldots, u_{r-1}$.

Sample principal components: Let $S$ denote the sample covariance matrix of the data matrix $X$ and let $Q=\left[q_{1}, \ldots, q_{p}\right]$ a $p \times p$ orthogonal matrix diagonalizing $S$ :

$$
Q^{T} S Q=\operatorname{diag}\left(\ell_{1}, \ldots, \ell_{p}\right)
$$

The $\ell_{j}$ are the sample variances that are estimates for $\lambda_{j}$. The vectors $q_{j}$ are sample estimates for the vectors $h_{j}$.
If $x$ is the random vector and $u=H^{T} x$ is the vector of principal components, then $\hat{u}=Q^{T} x$ is the vector of sample principal components.

## SCREE PLOTS



In applications: How many of the $\ell_{j}$ 's are significant?

## CANONICAL CORRELATION ANALYSIS (CCA)

H. Hotelling, 1936

Suppose a large data set is naturally decomposed into two groups. For example, $p \times 1$ random vectors $\vec{x}_{1}, \ldots, \vec{x}_{n}$ make up one set and $q \times 1$ random vectors $\vec{y}_{1}, \ldots, \vec{y}_{m}$ the other. We are interested in the correlations between these two data sets. For example, in medicine we might have $n$ measurements of age, height, and weight ( $p=3$ ) and $m$ measurements of systolic and diastolic blood pressures $(q=2)$. We are interested in what combination of the components of $x$ is most correlated with a combination of the components of $y$. Population Canonical Correlations: Let $x$ and $y$ be two random vectors of size $p \times 1$ and $q \times 1$, respectively. We assume $\mathbb{E}(x)=\mathbb{E}(y)=0$ and $p \leq q$. Form the $(p+q) \times 1$ vector

$$
\binom{x}{y}
$$

and its $(p+q) \times(p+q)$ covariance matrix

$$
\Sigma=\left(\begin{array}{ll}
\Sigma_{11} & \Sigma_{12} \\
\Sigma_{21} & \Sigma_{22}
\end{array}\right)
$$

Let

$$
u:=\alpha^{T} x \in \mathbb{R}, \quad v:=\gamma^{T} y \in \mathbb{R}
$$

where $\alpha$ and $\gamma$ are vectors to be determined. We want to maximize the correlation

$$
\operatorname{corr}(u, v)=\frac{\operatorname{cov}(u, v)}{\sqrt{\operatorname{var}(u) \operatorname{var}(v)}}
$$

The correlation does not change under scale transformations $u \rightarrow c u$, etc. so we can maximize this correlation subject to the constraints

$$
\begin{align*}
& \mathbb{E}\left(u^{2}\right)=\mathbb{E}\left(\alpha^{T} x \cdot \alpha^{T} x\right)=\alpha^{T} \Sigma_{11} \alpha=1  \tag{1}\\
& \mathbb{E}\left(v^{2}\right)=\gamma^{T} \Sigma_{22} \gamma=1 \tag{2}
\end{align*}
$$

Under these constraints

$$
\operatorname{corr}(u, v)=\mathbb{E}\left(\alpha^{T} x \cdot \gamma^{T} y\right)=\alpha^{T} \Sigma_{12} \gamma
$$

Let

$$
\psi=\alpha^{T} \Sigma_{12} \gamma-\frac{1}{2} \rho\left(\alpha^{T} \Sigma_{11} \alpha-1\right)-\frac{1}{2} \lambda\left(\gamma^{t} \Sigma_{22} \gamma-1\right)
$$

where $\rho$ and $\lambda$ are Lagrange multipliers. Set the vector of partial derivatives to zero:

$$
\begin{align*}
& \frac{\partial \psi}{\partial \alpha}=\Sigma_{12} \gamma-\rho \Sigma_{11} \alpha=0  \tag{3}\\
& \frac{\partial \psi}{\partial \gamma}=\Sigma_{12}^{T} \alpha-\lambda \Sigma_{22} \gamma=0 \tag{4}
\end{align*}
$$

If we left multiply (3) by $\alpha^{T}$ and (4) by $\gamma^{T}$, use the normalization conditions (1) and (2) we conclude $\lambda=\rho$. Thus (3) and (4) become

$$
\left(\begin{array}{rr}
-\rho \Sigma_{11} & \Sigma_{12}  \tag{5}\\
\Sigma_{21} & -\rho \Sigma_{22}
\end{array}\right)\binom{\alpha}{\gamma}=0
$$

with

$$
\operatorname{corr}(u, v)=\rho
$$

and $\rho \geq 0$ is a solution to

$$
\operatorname{det}\left(\begin{array}{rr}
-\rho \Sigma_{11} & \Sigma_{12} \\
\Sigma_{21} & -\rho \Sigma_{22}
\end{array}\right)=0
$$

This is a polynomial in $\rho$ of degree $(p+q)$. Let $\rho_{1}$ denote the maximum root and $\alpha_{1}$ and $\gamma_{1}$ corresponding solutions to (5).

Definition: $u_{1}$ and $v_{1}$ are called the first canonical variables and their correlation $\rho_{1}=\operatorname{corr}\left(u_{1}, v_{1}\right)$ is called the first canonical correlation coefficient.

More generally, the rth pair of canonical variables is the pair of linear combinations $u_{r}=\left(\alpha^{(r)}\right)^{T} x$ and $v_{r}=\left(\gamma^{(r)}\right)^{T} y$, each of unit variance and uncorrelated with the first $r-1$ pairs of canonical variables and having maximum correlation. The correlation $\operatorname{corr}\left(u_{r}, v_{r}\right)$ is the $r$ th canonical correlation coefficient.

## REDUCTION TO AN EIGENVALUE PROBLEM

Since

$$
\begin{gathered}
\left(\begin{array}{rr}
-\rho \Sigma_{11} & \Sigma_{12} \\
\Sigma_{21} & -\rho \Sigma_{22}
\end{array}\right)= \\
\left(\begin{array}{ll}
\Sigma_{11} & 0 \\
0 & 1
\end{array}\right)\left(\begin{array}{cc}
-\rho I & \Sigma_{11}^{-1} \Sigma_{12} \Sigma_{22}^{-1} \\
\Sigma_{21} & -\rho I
\end{array}\right)\left(\begin{array}{ll}
1 & 0 \\
0 & \Sigma_{22}
\end{array}\right)
\end{gathered}
$$

Thus the determinantal equation becomes

$$
\operatorname{det}\left(\rho^{2} I-\Sigma_{21} \Sigma_{11}^{-1} \Sigma_{12} \Sigma_{22}^{-1}\right)=0
$$

The nonzero roots $\rho_{1}, \ldots, \rho_{k}$ are called the population canonical correlation coefficients. $k=\operatorname{rank}\left(\Sigma_{12}\right)$.

In applications $\Sigma$ is not known. One uses the sample covariance matrix $S$ to obtain sample canonical correlation coefficients.

## AN EXAMPLE

The first application of Hotelling's canonical correlations is by F. Waugh [16] in 1942. He begins his paper with

Professor Hotelling's paper, "Relations between Two Sets of Variates," should be widely known and his method used by practical statisticians. Yet, few practical statisticians seem to know of the paper, and perhaps those few are inclined to regard it as a mathematical curiosity rather than an important and useful method for analyzing concrete problems. This may be due to ...

The Practical Problem: Relation of wheat characteristics to flour characteristics. Guiding principle

The grade of the raw material should give a good indication of the probable grade of the finished product.

The data: 138 samples of Canadian Hard Red Spring wheat and the flour made from each these samples. ${ }^{\text {a }}$

## wheat quality

$$
\begin{array}{lll}
x_{1}= & \text { kernel texture } & y_{1}=\text { wheat per barrels of flour } \\
x_{2}= & \text { test weight } & y_{2}=\text { ash in flour } \\
x_{3}= & \text { damaged kernels } & y_{3}=\text { crude protein in flour } \\
x_{4}= & \text { crude protein in wheat } & y_{4}=\text { gluten quality index } \\
x_{5}= & \text { foreign material } &
\end{array}
$$

$$
u_{1}=\alpha_{1}^{T} x=\text { index of wheat quality, } v_{1}=\gamma_{1}^{T} y=\text { index of flour quality }
$$

$$
\operatorname{corr}\left(u_{1}, v_{1}\right)=0.909
$$

[^0]
## DISTRIBUTION OF SAMPLE CANONICAL CORRELATION COEFFICIENTS

Give a sample of size $n$ observations on $\binom{x}{y}$ drawn from
$N_{p+q}(\mu, \Sigma)$ and $A$ the (unnormalized) sample covariance matrix.
Then $W$ is $W_{p+q}(n, \Sigma)$. We have [13]
Theorem (Constantine, 1963): Let $A$ have the $W_{p+q}(n, \Sigma)$
distribution where $p \leq q, n \geq p+q$ and $\Sigma$ and $A$ are partitioned as above. Then the joint probability density function of $r_{1}^{2}, \ldots, r_{p}^{2}$, the eigenvalues of $A_{11}^{-1} A_{12} A_{22}^{-1} A_{21}$ (let $\xi_{j}:=r_{j}^{2}$ ) is
$c_{p, q, n} \prod_{j=1}^{p}\left(1-\rho_{j}^{2}\right)^{n / 2} \prod_{j=1}^{p}\left[\xi_{j}^{(q-p-1) / 2}\left(1-\xi_{j}\right)^{(n-p-q-1) / 2}\right] \cdot|\Delta(\xi)| \cdot \mathcal{F}(\xi)$
where

- $c_{p, q, n}$ is a normalization constant.
- $\Delta(\xi)=$ Vandermonde determinant $=\prod_{j<k}^{p}\left(\xi_{j}-\xi_{k}\right)$.
- $\mathcal{F}(\xi)$ is a two-matrix hypergeometric function which can be expressed as an infinite series involving zonal polynomials. See Theorem 11.3.2 and Definition 7.3.2 in Muirhead [13].

Null Distribution: For $\Sigma_{12}=0$ ( $x$ and $y$ are independent), the above joint density for the sample canonical correlation coefficients reduces to

$$
c_{p, q, n} \prod_{j=1}^{p}\left[\xi_{j}^{(q-p-1) / 2}\left(1-\xi_{j}\right)^{(n-p-q-1) / 2}\right] \cdot|\Delta(\xi)|
$$

In this case the distribution of the largest sample canonical correlation coefficient $r_{1}$ can be used for testing the null hypothesis: $H: \Sigma_{12}=0$. We reject $H$ for large values of $r_{1}$.

## ZONAL POLYNOMIALS

Zonal polynomials naturally arise in multivariate analysis when considering group integrals such as

$$
\int_{\mathcal{O}(m)} \mathrm{e}^{-\operatorname{tr}\left(X H Y H^{T}\right)}(d H)
$$

where $X$ and $Y$ are $m \times m$ symmetric, positive definite matrices and $(d H)$ is normalized Haar measure. Zonal polynomials can be defined either through the representation theory of $G L(m, \mathbb{R})$ [12] or as eigenfunctions of certain Laplacians [5, 13].

Let $\mathcal{S}_{m}$ denote the space of $m \times m$ symmetric, positive definite matrices. Zonal polynomials, $C_{\lambda}(X), X \in \mathcal{S}_{m}$ are certain homogeneous polynomials in the eigenvalues of $X$ that are indexed by partitions $\lambda$.

To give the precise definition we need some preliminary definitions.

Definition: A partition $\lambda$ of $n$ is a sequence $\lambda=\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{\ell}\right)$ where the $\lambda_{j} \geq 0$ are weakly decreasing and $\sum_{j} \lambda_{j}=n$. We denote this by $\lambda \vdash n$.

For example, $\lambda=(5,3,3,1)$ is a partition of 12 . The number of nonzero parts of $\lambda$ is called the length of $\lambda$, denoted $\ell(\lambda)$.

If $\lambda$ and $\mu$ are two partitions of $n$, we say $\lambda<\mu$ (lexicographic order $)$ if, for some index $i, \lambda_{j}=\mu_{j}$ for $j<i$ and $\lambda_{j}<\mu_{j}$. For example

$$
(1,1,1,1)<(2,1,1)<(2,2)<(3,1)<(4)
$$

If $\lambda \vdash n$ with $\ell(\lambda)=m$, we define the monomial

$$
x^{\lambda}=x_{1}^{\lambda_{1}} x_{2}^{\lambda_{2}} \cdots x_{m}^{\lambda_{m}}
$$

and say $x^{\mu}$ is of higher weight than $x^{\lambda}$ if $\mu>\lambda$.

## Metrics and Laplacians

Let $X \in \mathcal{S}_{m}$ and $d X=\left(d x_{j k}\right)$ the matrix of differentials of $X$. We define a metric on $\mathcal{S}_{m}$ by

$$
(d s)^{2}=\operatorname{tr}\left(X^{-1} d X \cdot X^{-1} d X\right)
$$

A simple computation shows this metric is invariant under

$$
X \longrightarrow L X L^{T}, L \in G L(m, \mathbb{R}) .
$$

Let $n=m(m+1) / 2=\#$ of independent elements of $X$. We denote by $\operatorname{vec}(X) \in \mathbb{R}^{n}$ the column vector representation of $X$, e.g.

$$
X=\left(\begin{array}{ll}
x_{11} & x_{12} \\
x_{12} & x_{22}
\end{array}\right), \quad x:=\operatorname{vec}(X)=\left(\begin{array}{l}
x_{11} \\
x_{12} \\
x_{22}
\end{array}\right)
$$

The metric $(d s)^{2}$ is a quadratic differential in $d x$.

For example,

$$
(d s)^{2}=\left(\begin{array}{lll}
d x_{11} & d x_{12} & d x_{22}
\end{array}\right) \cdot G(x) \cdot\left(\begin{array}{c}
d x_{11} \\
d x_{12} \\
d x_{22}
\end{array}\right)
$$

where

$$
\begin{gathered}
G(x)=\left(g_{i j}\right)= \\
\frac{1}{\left(x_{11} x_{22}-x_{12}^{2}\right)^{2}}\left(\begin{array}{rcr}
x_{22}^{2} & -2 x_{12} x_{22} & x_{12}^{2} \\
-2 x_{12} x_{22} & 2 x_{11} x_{22}+2 x_{12}^{2} & -2 x_{11} x_{12} \\
x_{12}^{2} & -2 x_{11} x_{12} & x_{11}^{2}
\end{array}\right)
\end{gathered}
$$

Labeling $x=\operatorname{vec}(X)$ with a single index the differential is in the standard form

$$
(d s)^{2}=\sum_{i<j} d x_{i} g_{i j} d x_{j}
$$

and the Laplacian associated to this metric is

$$
\Delta_{X}=(\operatorname{det} G)^{-1 / 2} \sum_{j=1}^{n} \frac{\partial}{\partial x_{j}}\left[(\operatorname{det} G)^{1 / 2} \sum_{i=1}^{n} g^{i j} \frac{\partial}{\partial x_{i}}\right]
$$

where

$$
G^{-1}=\left(g^{i j}\right)
$$

More succinctly, if

$$
\begin{gathered}
\nabla_{X}=\left(\begin{array}{c}
\frac{\partial}{\partial x_{1}} \\
\vdots \\
\frac{\partial}{\partial x_{n}}
\end{array}\right) \\
\Delta_{X}=(\operatorname{det} G)^{-1 / 2}\left(\nabla_{X},(\operatorname{det} G)^{1 / 2} G^{-1} \nabla_{X}\right)
\end{gathered}
$$

where $(\cdot, \cdot)$ is the standard inner product on $\mathbb{R}^{n}$.

One can show that $\Delta_{X}$ is an invariant differential operator:

$$
\Delta_{L X L^{T}}=\Delta_{X}, \quad L \in G L(m, \mathbb{R})
$$

We now diagonalize $X$

$$
X=H Y H^{T}, \quad Y=\operatorname{diag}\left(y_{1}, \ldots, y_{m}\right), H \in \mathcal{O}(m) .
$$

The Laplacian $\Delta_{X}$ is now expressed in terms of a radial part and an angular part. The radial part of $\Delta_{X}$ is the differential operator

$$
\sum_{j=1}^{m} y_{j}^{2} \frac{\partial^{2}}{\partial y_{j}^{2}}+\sum_{j=1}^{m} \sum_{\substack{k=1 \\ k \neq j}}^{m} \frac{y_{j}^{2}}{y_{j}-y_{k}} \frac{\partial}{\partial y_{j}}+\sum_{j} y_{j} \frac{\partial}{\partial y_{j}}
$$

We now let $\Delta_{X}$ denote only the radial part.
We can now define zonal polynomials!

Definition: Let $X \in \mathcal{S}_{m}$ with eigenvalues $x_{1}, \ldots, x_{m}$ and let $\lambda=\left(\lambda_{1}, \ldots, \lambda-m\right) \vdash k$ into not more than $m$ parts. Then $C_{\lambda}(X)$ is the symmetric, homogeneous polynomial of degree $k$ in $x_{j}$ such that

1. The term of highest weight in $C_{\lambda}(X)$ is $x^{\lambda}$
2. $C_{\lambda}(X)$ is an eigenfunction of the Laplacian $\Delta_{X}$.
3. 

$$
(\operatorname{tr}(X))^{k}=\left(x_{1}+\cdots+x_{m}\right)^{k}=\sum_{\substack{\lambda \not-k \\ e(\lambda) \leq m}} C_{\lambda}(X) .
$$

## Remarks:

- Must show there is an unique polynomial satisfying these requirements.
- Eigenvalue in (2) equals $\alpha_{\lambda}:=\sum_{j} \lambda_{j}\left(\lambda_{j}-j\right)+k(m+1) / 2$.
- Program MOPS [5] computes zonal polynomials.

The following theorem is at the core of why zonal polynomials appear in multivariate statistical analysis.
Theorem: If $X, Y \in S_{p}$, then

$$
\begin{equation*}
\int_{\mathcal{O}(p)} C_{\lambda}\left(X H Y H^{T}\right)(d H)=\frac{C_{\lambda}(X) C_{\lambda}(Y)}{C_{\lambda}\left(I_{p}\right)} \tag{6}
\end{equation*}
$$

where $(d H)$ is normalized Haar measure.
Proof: Let $f_{\lambda}(Y)$ denote the left-hand side of (6) and $Q \in \mathcal{O}(p)$. $f_{\lambda}\left(Q Y Q^{T}\right)=f_{\lambda}(Y)$ (let $H \rightarrow H Q$ in integral and use invariance of the measure). Thus $f_{\lambda}$ is a symmetric function of the eigenvalues of $Y$. Since $C_{\lambda}$ is homogeneous of degree $|\lambda|$, so is $f_{\lambda}$. Now apply the

Laplacian to $f_{\lambda}$.

$$
\begin{aligned}
\Delta_{Y} f_{\lambda}(Y) & =\int_{\mathcal{O}(p)} \Delta_{Y} C_{\lambda}\left(X H Y H^{T}\right)(d H) \\
& =\int \Delta_{Y} C_{\lambda}\left(X^{1 / 2} H Y H^{T} X^{1 / 2}\right)(d H) \\
& =\int \Delta_{Y} C_{\lambda}\left(L Y L^{T}\right)(d H)\left(L=X^{1 / 2} H\right) \\
& =\int \Delta_{L Y L^{T}} C_{\lambda}\left(L Y L^{T}\right)(d H) \text { invariance of } \Delta_{Y} \\
& =\alpha_{\lambda} \int C_{\lambda}\left(L Y L^{T}\right)(d H) \\
& =\alpha_{\lambda} f_{\lambda}(Y)
\end{aligned}
$$

By definition of $C_{\lambda}(Y)$ we have $f_{\lambda}(Y)=d_{\lambda} C_{\lambda}(Y)$. Since $f_{\lambda}\left(I_{p}\right)=C_{\lambda}(X)$, we find $d_{\lambda}$ and the theorem follows.

Using Zonal Polynomials to Evaluate Group Integrals

$$
\begin{aligned}
\int_{\mathcal{O}(p)} & \mathrm{e}^{-\rho \operatorname{tr}\left(X H Y H^{T}\right)}(d H) \\
& =\sum_{k=0}^{\infty} \frac{\rho^{k}}{k!} \int_{\mathcal{O}(p)}\left(\operatorname{tr}\left(X H Y H^{t}\right)\right)^{k}(d H) \\
& =\sum_{k=0}^{\infty} \frac{\rho^{k}}{k!} \sum_{\substack{\lambda \vdash k \\
\ell(\lambda) \leq m}} \int_{\mathcal{O}(p)} C_{\lambda}\left(X H Y H^{T}\right)(d H) \\
= & \sum_{k=0}^{\infty} \frac{\rho^{k}}{k!} \sum_{\substack{\lambda \vdash k \\
\ell(\lambda) \leq m}} \frac{C_{\lambda}(X) C_{\lambda}(Y)}{C_{\lambda}\left(I_{p}\right)}
\end{aligned}
$$

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[^0]:    ${ }^{\text {a }}$ In this example $p>q$. The data are normalized to mean 0 and variance 1 .

