## EXPECTATION OF ELEMENTARY SYMMETRIC FUNCTIONS OF A WISHART MATRIX

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Some conjectures made by De Waal [Ann. Math. Statist. 43 (1972) 344-347] on the expectation of elementary symmetric functions of the roots of a noncentral Wishart matrix are proved true. The method of proof is through a simple, though perhaps obscure, property of these elementary symmetric functions and simple properties of the Wishart distribution.

1. Introduction. If  $X(p \times n)$  has independently, normally distributed columns with covariance V and  $\mathcal{E}(X) = M$ , then A = XX' has the noncentral Wishart distribution. Write  $\operatorname{tr}_j(A)$  for the jth elementary symmetric function of the latent roots of A and set  $\operatorname{tr}_0(A) = 1$  for convenience. Put  $\Omega = V^{-1}MM'$  and  $(a)^{(i)} = a(a-1)(a-2)\cdots(a-i+1)$ .

De Waal (1972) conjectures that

(1.1) 
$$\mathscr{E}\operatorname{tr}_{p}(A) = \mathscr{E}|A| = |V| \sum_{i=0}^{p} (n-i)^{(p-i)} \operatorname{tr}_{i}(\Omega), \quad p \leq n$$

and, further, that when V = I,

$$(1.2) \mathscr{E} \operatorname{tr}_{j}(A) = \sum_{i=0}^{j} (n-i)^{(j-i)} \binom{p-i}{j-i} \operatorname{tr}_{i}(\Omega), \quad i \leq j \leq p \leq n.$$

We will show

(1.3) 
$$\mathscr{E} \operatorname{tr}_{j}(V^{-1}A) = \sum_{i=0}^{j} (n-i)^{(j-i)} \binom{p-i}{j-i} \operatorname{tr}_{i}(\Omega), \quad i \leq j \leq p \leq n.$$

Labelling the claimed identity in (1.3) as  $C_{j,p}(V)$  it is clear that  $C_{p,p}(V)$  is equivalent to (1.1) and that  $C_{j,p}(I)$  is identical to (1.2). However, choosing K so that KVK' = I and transforming  $X \to KX$ , we see that  $C_{j,p}(V)$  is true if and only if  $C_{j,p}(I)$  is true  $1 \le j \le p$ . We show first that  $C_{p,p}(I)$  implies  $C_{j,p}(I)$  and then prove that  $C_{p,p}(I)$  is true thus validating the entire set of De Waal's conjectures, more generally expressed in (1.3).

2.  $C_{p,p}(I)$  implies  $C_{j,p}(I)$ . Let  $J=\{i_1,i_2,\cdots,i_j\}$  with  $i_1< i_2,\cdots,< i_j$  be an ordered subset of the integers  $1,2,\cdots,p$ . For any matrix  $B(p\times p)$ , define  $B_J$  as the  $j\times j$  matrix formed from B by preserving only those rows and columns corresponding to elements of J. If  $\lambda_1,\lambda_2,\cdots,\lambda_p$  are the latent roots of B then  $|B-\lambda I|=(\lambda_1-\lambda)(\lambda_2-\lambda)\cdots(\lambda_p-\lambda)$ . Differentiate this equation (p-j) times with respect to  $\lambda$  and set  $\lambda=0$  to obtain

$$(2.1) \operatorname{tr}_{i}(B) = \sum_{I} |B_{I}|$$

the summation extending over all possible J. From (2.1)

(2.2) 
$$\sum_{J} \operatorname{tr}_{k}(B_{J}) = \sum_{K(J)} |(B_{J})_{K(J)}|$$

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where K(J) is an ordered, size k subset of the elements of J. Clearly  $(B_J)_{K(J)} \equiv B_K$  for some K, an ordered, size k subset of  $\{1, 2, \dots, p\}$ . Further all possible such  $B_K$  will be included on the right-hand side of (2.2) an equal number of times. We have therefore

the combinational multiplier having been determined by comparing the number of determinants on the left and right sides of the equation. Use (2.1) on the right-hand side of (2.3) and then (2.2) to obtain

$$\sum_{J} \operatorname{tr}_{k}(B_{J}) = \binom{p-k}{j-k} \operatorname{tr}_{k}(B) \qquad k \leq j.$$

Assume  $C_{pp}(I)$  to be true,  $1 \le p \le n$ . Since

(2.5) 
$$\mathscr{E}\operatorname{tr}_{i}(A) = \sum_{J} \mathscr{E}|A_{J}|$$

then

(2.6) 
$$\mathscr{E}\operatorname{tr}_{i}(A) = \sum_{J} \sum_{i=0}^{J} (n-i)^{(j-i)} \operatorname{tr}_{i}(\Omega_{J})$$

and so,  $C_{j,n}(I)$  is true :  $1 \le j \le p \le n$  after using (2.4) on the right side of (2.6).

3. Validation of  $C_{p,p}(I)$ . Since  $\mathscr{C}|A|$  depends on  $\Omega$  only through the latent roots of the latter we may assume that  $\Omega=\operatorname{diag}\left(\omega_{1},\,\omega_{2},\,\cdots,\,\omega_{p}\right)$  and correspondingly that X=Y+M wherein Y is a matrix of pn mutually independent standard normal deviates and  $M=(m_{ij})$  is the matrix  $m_{ii}=\omega_{i}^{\frac{1}{2}},\ 1\leq i\leq p;$   $m_{ij}=0$  otherwise. Clearly  $\mathscr{C}|A|$  is a symmetric function of  $\{\omega_{1}^{\frac{1}{2}},\,\omega_{2}^{\frac{1}{2}},\,\cdots,\,\omega_{p}^{\frac{1}{2}}\}$  of order at most two in each element. Since the elements of Y have a distribution symmetric about zero,  $\mathscr{C}|A|$  is invariant under the transformation  $\omega_{i}^{\frac{1}{2}}\to -\omega_{i}^{\frac{1}{2}}$ . Evidently  $\mathscr{C}|A|$  must be a symmetric function of  $\{\omega_{1},\,\omega_{2},\,\cdots,\,\omega_{p}\}$  of order at most one in each element therefore, for some  $d_{0},\,d_{1},\,\cdots,\,d_{p}$ , we may write

(3.1) 
$$\mathscr{E}|A| = d_0 + d_1 \operatorname{tr}_1(\Omega) + \cdots + d_p \operatorname{tr}_p(\Omega).$$

We first show

(3.2) 
$$\mathscr{E} \frac{\partial}{\partial \omega_1} \frac{\partial}{\partial \omega_2}, \, \cdots, \, \frac{\partial}{\partial \omega_p} |A| = 1.$$

The differential of |A| with respect to  $\omega_1$  is the sum of p determinants the ith of which is |A| with its ith row differentiated with respect to  $\omega_1$ ,  $i=1,2,\cdots,p$ . The differential in (3.2) therefore is the sum of  $p^p$  determinants. Of these, those which have any row of A differentiated at least twice will be zero due to the presence of a row of zeros. Of the remaining p! determinants, all but one will have zero expectation since their expansion will be seen to contain an element of Y raised to an odd power. The surviving determinant is that in which the ith row of A has been differentiated with respect to  $\omega_i$ ,  $i=1,2,\cdots,p$ . The product of the diagonal elements of this determinant is  $(1+\omega_{11}^{\frac{1}{2}}y_{11})(1+\omega_{22}^{\frac{1}{2}}y_{22})\cdots (1+\omega_{2n}^{\frac{1}{2}}y_{nn})$  which has unit expectation. Every other product occurring in the

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expansion of this determinant has an element of Y raised to an odd power and has zero expectation.

Now

$$(3.3) d_i = \mathcal{E} \frac{\partial}{\partial \omega_1} \frac{\partial}{\partial \omega_2} \cdots \frac{\partial}{\partial \omega_i} |A| \text{when } \omega_{i+1} = \omega_{i+2} = \cdots = \omega_p = 0.$$

Let  $A_{11}$  be the first i rows of columns of A;  $A_{22}$  be the last p-i rows and columns of A and  $A_{12}=A'_{21}$  be formed from the first i rows and last p-i columns of A. When  $\omega_{i+1}=\omega_{i+2}=\cdots=\omega_p=0$ , it is well known that  $A_{11}$  and  $A_{22}-A_{21}A_{11}^{-1}A_{12}$  are independent and that the latter, which has a central Wishart distribution, has expectation  $(n-i)^{(p-i)}$ . Since  $|A|=|A_{11}||A_{22}-A_{21}A_{11}^{-1}A_{12}|$  and applying (3.1) to  $|A_{11}|$  we see that  $d_i=(n-i)^{(p-i)}$  which proves  $C_{pp}(I)$ .

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## REFERENCE

DE WAAL, D. J. (1972). On the expected values of the elementary symmetric functions of a noncentral Wishart matrix. *Ann. Math. Statist.* 43 344-347.

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