A Normal Limit Theorem For Moment Sequences

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ABSTRACT

Let Λ be the set of probability measures λ on [0,1]. Let $M_n = \{(c_1,\ldots,c_n) \mid \lambda \in \Lambda\}$ where $c_k = c_k(\lambda) = \int_0^1 x^k d\lambda$, $k = 1, 2, \ldots$ are the ordinary moments, and assign to the moment space M_n the uniform probability measure P_n . We show that, as $n \to \infty$, the fixed section (c_1,\ldots,c_k) properly normalized is asymptotically normally distributed. That is, $\sqrt{n}[(c_1,\ldots,c_k)-(c_1^0,\ldots,c_k^0)]$ converges to $\text{MVN}(0,\Sigma)$ where c_i^0 correspond to the arcsin law λ_0 on [0,1]. Properties of the $k \times k$ matrix Σ are given as well as some further discussion.

1. Introduction and Main Theorem

The set of probability measures on [0,1] is denoted as Λ , let further

$$M_n = \{(c_1, \dots, c_n) \mid \lambda \in \Lambda\}, \tag{1.1}$$

where $c_k = c_k(\lambda) = \int_0^1 x^k \lambda(dx)$, k = 0, 1, 2, ...; $(c_0 = 1)$. This so-called moment space M_n is the convex hull of the curve $\{(x, x^2, ..., x^n): 0 \le x \le 1\}$ in \mathbb{R}_n and is a very small compact subset of the unit cube $[0, 1]^n$. For instance, it is known that

$$V_n = \text{Vol } M_n = \prod_{k=1}^n B(k,k) = \prod_{k=1}^n \frac{\Gamma(k)\Gamma(k)}{\Gamma(2k)},$$
 (1.2)

see Karlin and Studden 1966, p. 129, Theorem 6.2; (another proof is given below). Thus V_n is roughly of size 2^{-n^2} , more precisely, $\log V_n \approx -n^2 \log 2$ as $n \to \infty$.

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Our investigations stem from an attempt to understand more fully the shape and structure of M_n by looking, in some sense, at a typical point of M_n . Let P_n be the uniform probability measure on M_n , i.e., $dP_n = dx/V_n$ is n-dimensional Lebesgue measure on M_n normalized by the volume of M_n . In this way $(c_1, \ldots, c_n) \in M_n$ can now be viewed as a random vector. The symbol E_n will indicate expected values relative to P_n .

For example, M_2 is determined by the inequalities $c_1^2 \le c_2 \le c_1 \le 1$ and has volume $V_2 = \frac{1}{6}$ thus $dP_2 = 6dc_1dc_2$ on M_2 . The marginal densities of c_1 , c_2 are $6(c_1 - c_1^2)$, $0 < c_1 < 1$, and $6(\sqrt{c_2} - c_2)$, $0 < c_2 < 1$, respectively. The means are $E_2[c_1] = 1/2$ and $E_2[c_2] = 2/5$ and the squared correlation is 35/38. General closed form expressions even for, say, the means $E_n[c_k]$ seem difficult to obtain.

The so-called center (c_1^0, \ldots, c_n^0) of the moment space M_n is given by

$$c_k^0 = \int_0^1 x^k f_0(x) dx = 2^{-2k} {2k \choose k} \approx \frac{1}{\sqrt{\pi k}} \text{ as } k \to \infty.$$
 (1.3)

Here, $f_0(x) = \pi^{-1}x^{-1/2}(1-x)^{-1/2}$, 0 < x < 1, is the density of the arc-sin probability measure λ_0 on [0,1]. The word "center" will become clearer below. Our main result is the following.

Theorem 1.1. As $n \to \infty$, the distribution of $\sqrt{n}[(c_1, \ldots, c_k) - (c_1^0, \ldots, c_k^0)]$ relative to P_n converges to a multivariate normal distribution $MVN(0, \Sigma_k)$. Here, $\Sigma_k = \frac{1}{2}A_kA'_k$ with A_k as the lower triangular $k \times k$ matrix defined by

$$a_{ij} = 2^{-2i+1} {2i \choose i-j} \text{ if } 1 \le j \le i;$$

$$= 0 \qquad \text{if } j > i;$$

$$(1.4)$$

thus $a_{ii}=2^{-2i+1}$. In particular, if c_k is governed by P_n and $n\to\infty$ then $c_k\to c_k^0$ in probability.

By $A=(a_{ij};\ 1\leq i,\ j<\infty)$ we will denote the corresponding infinite lower triangular matrix, having A_k as its left upper $k\times k$ submatrix. The proof of the theorem is, in essence, quite simple and, at the same time, illuminating. The boundary of M_n has P_n -measure zero and thus can be ignored. Note that $(c_1,\ldots,c_n)\in \operatorname{int} M_n$ implies that $(c_1,\ldots,c_k)\in \operatorname{int} M_k$ for all $k\leq n$.

It will be convenient to employ the canonical coordinates $p_k(k=1,2,\ldots)$ introduced by Skibinsky (1967). For each $k=1,\ldots,n$, the k-th canonical coordinate p_k of a moment point $(c_1,\ldots,c_n)\in \operatorname{int} M_n$ is well-defined, satisfies $0< p_k<1$, and depends only on c_1,\ldots,c_k . The associated function $p_k=f_k(c_1,\ldots,c_k)$ is independent of n. Conversely, c_k is fully determined by p_1,\ldots,p_k .

Given $(c_1,\ldots,c_{k-1})\in M_{k-1}$, let $c_k^+=c_k^+(c_1,\ldots,c_{k-1})$ and $c_k^-=c_k^-(c_1,\ldots,c_{k-1})$, respectively, denote the largest and smallest possible value of c_k which is compatible with $(c_1,\ldots,c_{k-1},c_k)\in M_k$. Thus, $c_k^-\leq c_k\leq c_k^+$ when $(c_1,\ldots,c_k)\in M_k$. In particular, $c_1^-=0$; $c_1^+=1$ and $c_2^-=c_1^2$; $c_2^+=c_1$. As is easily seen, $(c_1,\ldots,c_k)\in \mathrm{int}\,M_k$ if and only if $c_j^-< c_j< c_j^+(j=1,\ldots,k)$. Put

$$\Delta_k = \Delta_k(c_1, \ldots, c_{k-1}) = c_k^+(c_1, \ldots, c_{k-1}) - c_k^-(c_1, \ldots, c_{k-1}).$$

Here, $\Delta_k > 0$ for all $(c_1, \ldots, c_{k-1}) \in \text{ int } M_{k-1}$. For $k = 1, \ldots, n$, the k-th canonical coordinate (or moment) of a moment point $(c_1, \ldots, c_n) \in \text{ int } M_n$ is defined by

$$p_k = \frac{c_k - c_k^-}{c_k^+ - c_k^-} \text{ thus } c_k = c_k^-(c_1, \dots, c_{k-1}) + \Delta_k(c_1, \dots, c_{k-1}) p_k.$$
 (1.5)

Note that $0 < p_k < 1$. It follows by induction that, for all $k \ge 1$, there is a 1:1 correspondence between points $(c_1, \ldots, c_k) \in \text{int } M_k$ and points $(p_1, \ldots, p_k) \in (0, 1)^k$. Thus c_k^- , c_k^+ and $\Delta_k = c_k^+ - c_k^-$ can also be regarded as functions of p_1, \ldots, p_{k-1} ; these functions happen to be polynomial, (as is clear from (3.6) or (3.19)). Similarly, c_k is a polynomial in p_1, \ldots, p_k which is linear in the variable p_k with coefficient Δ_k , (see (1.5)). The canonical moments p_k for the Beta (α, β) distribution on [0,1] are given in Skibinsky (1969) p. 1759. The above arc—sin distribution λ_0 corresponds to $\alpha = \beta = 1/2$ and has canonical moments $p_k^0 = 1/2$ for all $k \ge 1$. This partially explains why the corresponding moment point (c_1^0, \ldots, c_n^0) may be regarded to be the center of M_n . Here, the c_k^0 are as in (1.3).

Remark. The canonical coordinates p_k admit a more general interpretation and as such are quite robust. Namely, consider any non-degenerate compact interval [a,b] and let $\{W_j(x)\}_{j=1}^{\infty}$ be a given system of polynomials of the form $W_j(x) = \sum_{m=0}^{j} d_{jm}x^m$ with $d_{jj} > 0$. For example, $W_j(x) = x^j$. Next consider all moment sequences $\{w_j\}_{j=1}^{\infty}$ of the

form $w_j = \int W_j(x)\lambda(dx)$ $(j=1,2,\ldots)$ with λ as a probability measure on [a,b]. Given the moments w_1,\ldots,w_{n-1} , let w_n^- , w_n^+ denote the smallest and largest possible value of w_n . Provided $\Delta_n = w_n^+ - w_n^- > 0$, define $p_n = (w_n - w_n^-)/\Delta_n$; thus $0 \le p_n \le 1$. As is easily seen, the resulting sequence $\{p_n\}$ of (generalized) canonical coordinates is independent of the particular choice of the system of polynomials $\{W_j(x)\}$. In addition, as was already observed by Skibinsky (1969) p.1763 Theorem 5, if the probability measure λ on [a,b] is linearly transformed (with positive slope) to a measure μ on another interval $[\alpha, \beta]$ then λ and μ have exactly the same canonical coordinates $p_n(n \ge 1)$. Here, $\mu(F) = \lambda(g^{-1}F)$ where $g(x) = \alpha + (\beta - \alpha)(x - a)/(b - a)$.

Let us return to the above (Hausdorff) sequences $\{c_n\}$ of the special form $c_n = \int x^n \lambda(dx)$, with λ as a probability measure on [0,1]. Using (1.5), one finds that

$$\frac{\partial c_k}{\partial p_j} = 0 if j > k;$$

$$= \Delta_k = c_k^+ - c_k^- = \prod_{r=1}^{k-1} p_r q_r if j = k;$$
(1.6)

Here and from now on, $q_r = 1 - p_r$. The latter elegant formula for Δ_k was established by Skibinsky (1967). A different proof is given below, see (3.4). It follows from (1.6) that

$$\frac{\partial(c_1,\ldots,c_n)}{\partial(p_1,\ldots,p_n)} = \prod_{k=1}^n \frac{\partial c_k}{\partial p_k} = \prod_{r=1}^{n-1} (p_r q_r)^{n-r}.$$
 (1.7)

Transforming the integral $V_n = \int_{M_n} dc_1 \dots dc_n$ to an integral over $(0,1)^n$ relative to the p_j , we see that formula (1.2) above is an immediate consequence of (1.7). Both (1.2) and (1.7) are special cases of the following result, (namely, with m = 0 and m = n - 1, respectively).

Theorem 1.2. Let $0 \le m < n$ and $(c_1, \ldots, c_m) \in \text{ int } M_m$. Then the set $M_n(c_1, \ldots, c_m)$ of all (c_{m+1}, \ldots, c_n) such that $(c_1, \ldots, c_n) \in M_n$ has (n-m)-dimensional volume

Vol
$$M_n(c_1, ..., c_m) = \prod_{r=1}^m (p_r q_r)^{n-m} \prod_{k=2}^{n-m} \frac{\Gamma(k)\Gamma(k)}{\Gamma(2k)}.$$
 (1.8)

The latter is maximal when $p_r = 1/2$ (r = 1, ..., m). Note that, under P_n the conditional distribution of $(c_{m+1}, ..., c_n)$ given $(c_1, ..., c_m)$ is the uniform distribution $dc_{m+1} ... dc_n/\text{Vol}$ $M_n(c_1, ..., c_m)$ on $M_n(c_1, ..., c_m)$.

In the sequel, for each fixed n, when we assign to M_n the uniform distribution P_n , functions on M_n such as c_1, \ldots, c_k or p_1, \ldots, p_k $(k \leq n)$ can be regarded as random variables. But note that the resulting joint distribution will depend on n.

Proof. Prescribing $(c_1, \ldots, c_m) \in \text{ int } M_m$ is the same as prescribing the parameters $0 < p_r < 1 (r = 1, \ldots, m)$. Further note, using (1.6), that

$$\frac{\partial(c_{m+1},\ldots,c_n)}{\partial(p_{m+1},\ldots,p_n)} = \prod_{s=m+1}^n \prod_{r=1}^{s-1} p_r q_r = \prod_{r=1}^m (p_r q_r)^{n-m} \prod_{r=m+1}^{n-1} (p_r q_r)^{n-r}.$$

The volume on hand is equal to the integral of $dc_{m+1} \dots dc_n$ over $M_n(c_1, \dots, c_m)$. Transforming that integral to an integral with respect to the variables p_{m+1}, \dots, p_n over the unit cube $(0,1)^{n-m}$, one obtains (1.8).

Theorem 1.3. The uniform probability measure P_n on M_n is equivalent to the first n canonical coordinates p_1, \ldots, p_n being independent random variables in such a way that p_k has a symmetric Beta (α_k, α_k) distribution with $\alpha_k = n - k + 1, \ k = 1, \ldots, n$.

Proof. Simply transform the integral

$$E_n f(p_1, \ldots, p_n) = \int_{M_n} f(p_1, \ldots, p_n) dc_1 \ldots dc_n / V_n$$

where f is arbitrary, to the variables p_1, \ldots, p_n , again using (1.7).

The symmetric distribution Beta (α, α) $(\alpha > 0)$ has mean 1/2 and variance 1/(8 α +4). Hence, for k = 1, ..., n, letting $\alpha = n - k + 1$,

$$E_n[p_k] = \frac{1}{2}; \ \operatorname{Var}[p_k] = \frac{1}{8(n-k+3/2)} = \frac{1}{8n} + O\left(\frac{1}{n^2}\right),$$
 (1.9)

as $n \to \infty$. Moreover, as is well known and easily seen, $\sqrt{n} [p_k - 1/2] \to N(0, 1/8)$ in distribution, under P_n as $n \to \infty$. Two proofs of the following central Lemma are given in Section 3.

Lemma 1.4. The first order Taylor expansion of $c_k = c_k(p_1, \ldots, p_k)$ about the center (p_1^0, \ldots, p_k^0) with $p_j^0 = 1/2$ is given by

$$c_k = c_k^0 + 2\sum_{m=1}^k a_{km}(p_m - \frac{1}{2}) + O\left(\sum_{m=1}^k |p_m - \frac{1}{2}|^2\right).$$
 (1.10)

Here, the a_{km} are as in (1.4). In particular $a_{km} = 2^{-2k+1} {2k \choose k-m}$ if $m \le k$.

Proof of Theorem 1.1. Let k be fixed and j, m = 1, ..., k. With $n \ge k$ and relative to P_n as the underlying measure, consider the random variables $X_{nj} = \sqrt{n} (c_j - c_j^0)$ and $Z_{nm} = 2\sqrt{n}(p_m - 1/2)$. Here, $Z_{n1}, ..., Z_{nk}$ are independent, for each fixed n, while $Z_{nm} \to N(0, 1/2)$ when m is fixed and $n \to \infty$. Writing (1.10) as

$$X_{nj} = \sum_{m=1}^{k} a_{jm} Z_{nm} + O\left(\frac{1}{\sqrt{n}} \sum_{m=1}^{k} Z_{nm}^{2}\right), \qquad (j = 1, \dots, k),$$

Theorem 1.1 becomes an immediate consequence.

2. Further Discussion

Let Σ be the infinite symmetric matrix $\Sigma = (\sigma_{ij}) = \frac{1}{2}AA'$ having $\Sigma_k = \frac{1}{2}A_kA'_k$ as its left upper $k \times k$ submatrix. Recall that Σ_k is the covariance matrix of the asymptotic $\text{MVN}(0, \Sigma_k)$ distribution as $n \to \infty$ of $\sqrt{n} [(c_1, \ldots, c_k) - (c_1^0, \ldots, c_k^0)]$, when the latter is governed by the uniform measure P_n on M_n . Thus asymptotically, as $n \to \infty$, the c_i have means $c_i^0 + o(1)$ and covariances $(\sigma_{ij}/n)(1 + o(1))$. Let further

$$\rho_{ij} = \sigma_{ij} / \sqrt{\sigma_{ii}\sigma_{jj}}.$$

Thus ρ_{ij} is the limiting value as $n \to \infty$ of the correlation coefficient under P_n between the moments c_i and c_j . The following result is proved in Section 4.

Lemma 2.1. One has

$$\sigma_{ij} = c_{i+j}^0 - c_i^0 c_j^0, \tag{2.1}$$

where the c_k^0 are as in (1.3). Hence, $\sigma_{ij} \to 0$ as $i, j \to \infty$. If s is fixed then $\rho_{s,s+r} \to 0$ as $r \to \infty$. If r is fixed then $\rho_{s,s+r} \to 1$ as $s \to \infty$. More generally, for any fixed $\theta \in [0,1]$,

$$\rho_{ij} \to \left(\frac{4K}{(K+1)^2}\right)^{1/4} \text{ when } i, j \to \infty; \ j/i \to K.$$
(2.2)

Let $k \geq 1$ be fixed. It is natural to inquire into the diagonalization of the symmetric $k \times k$ matrix Σ_k and corresponding linear transformations of (c_1, \ldots, c_k) . In view of the

usual Gram-Schmidt orthogonalization procedure, it suffices to determine the essentially unique linear combinations $t_i = b_{i1}c_1 + \ldots + b_{ii}c_i$ $(1 \le i \le k)$ with $b_{ii} \ne 0$ that are asymptotically uncorrelated, under P_n as $n \to \infty$. Equivalently, letting $b_{im} = 0$ when m > i, we want $B_k = (b_{im}; i, m = 1, \ldots, k)$ to be a non-singular lower triangular $k \times k$ matrix such that $D_k = B_k \Sigma_k B'_k$ is diagonal. Adding suitable constants b_{io} , one can further achieve that

$$t_i = \sum_{m=0}^{i} b_{im} c_m \quad (i = 1, \dots, k; \ c_0 = 1)$$
 (2.3)

are asymptotically uncorrelated and of mean 0. Equivalently, letting $t_0 = c_0 = 1$, we want t_0, t_1, \ldots, t_k to be asymptotically orthogonal, under P_n as $n \to \infty$.

The above diagonalization process happens to be intimately connected with the usual Chebyshev polynomials. Namely, consider the probability space Ω_0 consisting of the interval [0, 1] together with the arc-sin measure λ_0 as the underlying probability measure. The functions $x \to x^i$ on Ω_0 can then be regarded as random variables Z_i . We see from (1.3) that $EZ_i = c_i^0$ and $EZ_iZ_j = c_{i+j}^0$. Therefore,

$$Cov(Z_i, Z_i) = \sigma_{ij} \text{ for all } i, j \ge 1, \tag{2.4}$$

with σ_{ij} exactly as in (2.1). Hence, the means and covariances, of $\sqrt{n}(c_i - c_i^0)(i = 1, ..., k)$ under P_n , coincide asymptotically (as $n \to \infty$) with the means and covariances of $Z_i - c_i^0$ (i = 1, ..., k). Thus the above diagonalization is equivalent to finding k + 1 linear combinations of the form $T_i^{\#} = \sum_{m=0}^{i} b_{im} Z_m$ (i = 0, 1, ..., k), with $b_{ii} \neq 0$; $b_{00} = 1$, that are orthogonal as random variables on Ω_0 . But that simply means that the corresponding polynomials

$$T_i^*(x) = \sum_{m=0}^i b_{im} x^m, \ (i = 0, 1, 2, ...),$$
 (2.5)

one of each degree, are orthogonal with respect to the arc-sin measure λ_0 . Choosing the leading coefficient b_{ii} appropriately, we may as well assume that the $T_i^*(x)$ are precisely the Chebyshev polynomials, adapted to the interval [0, 1]. And then the resulting coefficients b_{im} are independent of k, (where $k \geq \max(i, m)$).

The functions $\cos i\theta$ (i=0,1,2,...) are clearly orthogonal with respect to the uniform measure on $[0,\pi]$. Letting $y=\cos\theta$, $\cos i\theta=T_i(y)$ one arrives at the system $\{T_i(y)\}_{i=0}^{\infty}$

of ordinary Chebyshev polynomials, orthogonal with respect to the measure $dy/\sqrt{1-y^2}$ on (-1,1). Letting $x=(1+y)/2=(1+\cos\theta)/2=(\cos\theta/2)^2$, leads to the desired system

$$T_i^*(x) = T_i(2x-1), (i=0,1,\ldots)$$
 (2.6)

as in (2.5) of orthogonal polynomials with respect to the measure λ_0 on (0, 1). Here, $T_i^*(x)$ is of exact degree i, while $T_0^*(x) \equiv 1$. The coefficients in (2.5) are given by $b_{i0} = (-1)^i$ and

$$b_{im} = (-1)^{i+m} \ 2^{2m-1} \frac{i}{m} \binom{i+m-1}{i-m}$$

$$= (-1)^{i+m} \ 2^{2m} \frac{i}{i+m} \binom{i+m}{i-m} \text{ if } 1 \le m \le i.$$
(2.7)

Thus $b_{ii} = 2^{2i-1}$ if $i \ge 1$. Further, from now on, $b_{im} = 0$ if m > i. Formula (2.7) easily follows from the known result that $T_n(2x-1) = (-1)^n F(-n, n; \frac{1}{2}, x)$, see Abramowitz and Stegun (1965) p. 795 and Henrici (1977) p. 176. For the sake of completeness, an independent proof of (2.7) is included in Section 4. Further note that

$$\int_0^1 T_j^*(x)^2 \lambda_0(dx) = \int_0^\pi (\cos j\theta)^2 \frac{d\theta}{\pi} = \frac{1}{2}.$$
 (2.8)

Theorem 2.2. Consider the linear combinations

$$t_i = \sum_{m=0}^{i} b_{im} c_m = \sum_{m=1}^{i} b_{im} (c_m - c_m^0), \ (i = 1, 2, \dots; c_0 = 1).$$
 (2.9)

Here the b_{im} are as in (2.5) and (2.7). Then, for any fixed $k \geq 1$ and $n \to \infty$, the distribution of \sqrt{n} (t_1, \ldots, t_k) relative to P_n converges in distribution to the multivariate normal distribution $MVN(0, \frac{1}{2}I_k)$. Here, I_k denotes the $k \times k$ identity matrix.

Proof. The second equality sign in (2.9) follows from $c_0 = c_0^0 = 1$ and

$$t_{i}^{0} = \sum_{m=0}^{i} b_{im} c_{m}^{0} = \int_{0}^{1} T_{i}^{*}(x) \lambda_{0}(dx) = 0 \text{ if } i \ge 1.$$
 (2.10)

In view of Theorem 1.1, it suffices to show that $B_k \Sigma_k B'_k = \frac{1}{2} I_k$. In some sense this already follows from the previous discussion. As a direct proof, if $1 \le i, j \le k$ then

$$\sum_{r=0}^{k} \sum_{s=0}^{k} b_{ir} b_{js} (c_{i+j}^{0} - c_{i}^{0} c_{j}^{0}) = \sum_{r=0}^{k} \sum_{s=0}^{k} b_{ir} b_{js} c_{i+j}^{0} = \int_{0}^{1} T_{i}^{*}(x) T_{j}^{*}(x) \lambda_{0}(dx) = \frac{1}{2} \delta_{i}^{j}.$$

9

Here, we used (2.5), (2.8), (2.10) as well as the orthogonality of the $T_i^*(x)$ with respect to λ_0 . Note that $c_{i+j}^0 - c_i^0 c_j^0 = 0$ when either i = 0 or j = 0. In view of (2.1), it follows that $B_k \Sigma_k B_k' = \frac{1}{2} I_k$.

Theorem 2.3. The lower triangular matrices $A = (a_{ij}; i, j \ge 1)$ and $B = (b_{ij}; i, j \ge 1)$ are each others inverse. Similarly for A_k and B_k , (any $k \ge 1$). Moreover, for $m \ge 1$,

$$x^{m} = c_{m}^{0} + \sum_{r=1}^{m} a_{mr} T_{r}^{*}(x);$$
(2.11)

COROLLARY 2.4. We have for all $m, r \ge 1$ that

$$\int_0^1 x^m T_r^*(x) \lambda_0(dx) = \frac{1}{2} a_{mr}. \tag{2.12}$$

Moreover,

$$c_m = c_m^0 + \sum_{r=1}^m a_{mr} t_r. (2.13)$$

Here, the t_r are as in (2.3) thus $t_r = \int T_r^*(x) \lambda(dx)$.

We will present several proofs. Note that (2.12) is an immediate consequence of (2.8), (2.11) and the orthogonality of the $T_r^*(x)$ with respect to λ_0 . Further, (2.13) follows from (2.11) from an integration relative to any $\lambda \in \Lambda$ having the moments $c_0 = 1, c_1, \ldots, c_m$. Choosing $\lambda = \lambda_0$, one has $c_m = c_m^0(m \geq 0)$ and $t_r = 0$, $(r \geq 1; t_0 = 1)$. This explains the constant term c_m^0 in (2.11), (2.13). Finally observe that (2.11) is actually equivalent to Λ , B being each others inverse, as can be seen by substituting formula (2.5) for the $T_r^*(x)$ into (2.11), and equating coefficients.

A first proof of Theorem 2.3 amounts to a direct verification of (2.11), see Section 4. A second proof is to directly verify the property AB = I, see Section 4. As still another demonstration, recall that, in the above proof of Theorem 2.2, we already established that $B\Sigma B' = \frac{1}{2}I$ where $\Sigma = \frac{1}{2}AA'$. Hence, the lower triangular matrix C = BA satisfies CC' = I, in particular, the rows of C are mutually orthogonal. Also using that $c_{ii} = a_{ii}b_{ii} = (2^{-2i+1})(2^{2i-1}) = 1$, we conclude that C must be the identity matrix.

10

3. Proof of Lemma 1.4.

We will present two different proofs. The first one exploits an important relation between the Hausdorff moment problem and a certain random walk. This relation, which one of us plans to discuss in more detail in a subsequent paper, is implicit in the work of Karlin and McGregor (1959).

Let $\{X_n\}_{n=0}^{\infty}$ be a stationary discrete time Markov chain (also called random walk) on the nonnegative integers Z_+ which is determined by the transition probabilities

$$P(X_{n+1} = j \mid X_n = i) = p_i \quad \text{if } j = i - 1;$$

$$= q_i \quad \text{if } j = i + 1;$$

$$= 0, \quad \text{otherwise.}$$

$$(3.1)$$

Here, $q_i = 1 - p_i$. Further $0 < p_i < 1$ for $i \ge 1$, while $p_0 = 0$; $q_0 = 1$. The corresponding n-step probabilities are denoted as $P_{ij}^{(n)} = P(X_n = j \mid X_0 = i)$. It was shown by Karlin and McGregor (1959) p. 69 that there exists a necessarily unique probability measure λ of infinite support on [0,1] such that

$$P_{00}^{(2n)} = P(X_{2n} = 0 \mid X_0 = 0) = \int_0^1 x^n \lambda(dx), \text{ for all } n \ge 0.$$
 (3.2)

In other words,

$$c_n = P_{00}^{(2n)} \quad (n = 0, 1, \ldots)$$
 (3.3)

always defines a Hausdorff moment sequence having $c_0 = 1$; $(c_1, \ldots, c_n) \in \text{ int } M_n$ for all $n \geq 1$. In fact, (3.2) establishes a 1:1 correspondence between all such Hausdorff moment sequences $\{c_n\}$ on the one hand and all random walks $\{X_n\}$ on the other hand, each random walk being determined as above by a sequence $\{p_n\}_{n=1}^{\infty}$ of canonical coordinates, $0 < p_n < 1$.

Consider a random walk $\{X_n\}$ as above and define c_n as in (3.3). Conditional on $X_0 = 0$, the conditional probability $c_k = P_{00}^{(2k)}$, (to be back in state 0 after 2k steps, not necessarily for the first time), obviously depends only on the parameters p_1, \ldots, p_k . Fixing c_1, \ldots, c_k is equivalent to fixing p_1, \ldots, p_k . Hence, for given c_1, \ldots, c_{n-1} , the smallest and largest possible value c_n^- and c_n^+ of $c_n = P_{00}^{(2n)}$ is realized by choosing $p_n = 0$ or $p_n = 1$,

respectively. In fact, c_n^- represents the (common) part of the return probability $c_n = P_{00}^{(2n)}$ arising from paths of length 2n (from 0 back to 0 in 2n steps) which never reach state n, and thus have their probability as a function of p_1, \ldots, p_{n-1} , independent of p_n . Similarly, $c_n^+ - c_n^-$ is equal to the probability $q_1 q_2 \ldots q_{n-1} p_n p_{n-1} \ldots p_1$ of the *single* path which leads from 0 to 0 in 2n steps which does reach state n. Maximizing c_n given p_1, \ldots, p_{n-1} , that is, choosing $p_n = 1$, this reduces to

$$\Delta_n := c_n^+ - c_n^- = q_1 q_2 \dots q_{n-1} p_{n-1} p_{n-2} \dots p_1 = \prod_{r=1}^{n-1} p_r q_r > 0.$$
 (3.4)

Finally note that $c_n = P_{00}^{(2n)} = c_n^- + p_n(c_n^+ - c_n^-)$. Comparing the latter with (1.5), we conclude that, for all $n \geq 1$, the random walk parameter p_n coincides with the n-th canonical coordinate of the moment point $(c_1, \ldots, c_n) \in \text{int } (M_n)$.

First proof of Lemma 1.4. Let $\{c_n\}_{n=0}^{\infty}$ be a Hausdorff moment sequence and $\{p_r\}_{r=1}^{\infty}$ be the associated sequence of canonical coordinates. Let $r \geq 1$ be fixed and

$$C_n(r) = \left[\frac{\partial}{\partial p_r} c_n\right]_0 = \left[\frac{\partial}{\partial p_r} P_{00}^{(2n)}\right]_0. \tag{3.5}$$

The subscript zero here indicates that $p_k = p_k^0 = 1/2$, for all $k \ge 1$. We want to show that $C_n(r) = 2a_{nr}$ with a_{nr} as in (1.4). In the present proof, we exploit the above random walk interpretation. Hence,

$$c_n = P_{00}^{(2n)} = \sum_{n=0}^{\infty} p_1^{m_1} p_2^{m_2} \dots q_0^{n_0} q_1^{n_1} \dots,$$
 (3.6)

where we sum over all paths $x=(x_0,x_1,\ldots,x_{2n})$ with $x_k-x_{k-1}=\pm 1(k=1,\ldots,2n)$ and such that $x_0=0$; $x_{2n}=0$; (thus c_n is a polynomial of degree 2n-1 in terms of p_1,\ldots,p_n). Further, for each such path, $m_j(j\geq 1)$ and $n_j(j\geq 0)$, respectively, will denote the number of transitions $x_{k-1}\to x_k(k=1,\ldots,2n)$ of type $j\to j-1$ and $j\to j+1$, respectively. Differentiating the latter sum with respect to p_r causes an extra factor $\frac{m_r}{p_r}-\frac{n_r}{q_r}$. ((tting afterwards $p_k=\frac{1}{2}$ for all $k\geq 1$, we find that

$$C_n(r) = 2E((m_r - n_r)I_0(X_{2n}) \mid X_0 = 0).$$
(3.7)

where $I_0(x)$ is the indicator function on the set $\{0\}$. Here, and from now on in the present proof, $\{X_n\}$ will be the simple random walk on Z_+ having 1-step probabilities $p_k = q_k = \frac{1}{2}$

for all $k \geq 1$, (while $p_0 = 0$; $q_0 = 1$). Moreover, since the path $\{X_0, X_1, \ldots, X_{2n}\}$ is random so are the associated transition numbers m_j and n_j .

Let further $\{Y_n\}_{n=0}^{\infty}$ be the classical random walk on $Z = \{0, \pm 1, \pm 2, ...\}$ with independent increments such that $P(Y_n - Y_{n-1} = -1) = P(Y_n - Y_{n-1} = +1) = \frac{1}{2}$. For each $s \in Z$, let

$$D_n(s) = E[(m_s - n_s)I_0(Y_{2n}) \mid Y_0 = 0].$$
(3.8)

Here, m_s and n_s , respectively, denote the (random) number of transitions $Y_{k-1} \to Y_k (k = 1, ..., 2n)$ of the form $s \to s-1$ and $s \to s+1$, respectively.

Identifying the states j and -j (for all j), the process $\{Y_n\}$ reduces precisely to the above simple random walk $\{X_n\}$. And it easily follows from (3.7) that

$$\frac{1}{2}C_n(r) = D_n(r) - D_n(-r) = 2D_n(r). \tag{3.9}$$

We further claim that

$$D_n(r) = P(Y_{2n} = 0; Y_k = r \text{ for some } 0 \le k \le 2n \mid Y_0 = 0).$$
 (3.10)

After all, consider any fixed path $y = (y_0, y_1, \ldots, y_{2n})$ with $y_k - y_{k-1} = \pm 1 (k = 1, \ldots, 2n)$ and $y_0 = 0$; $y_{2n} = 0$. Since $r \ge 1$ such a path y can contribute to $D_n(r)$ only when $y_k = r$ for some $0 \le k \le n$. Let k_1 and k_2 be the minimal and maximal such index k. Thus, $0 < k_1 \le k_2 < 2n$ and further $y_{k_1} = y_{k_2} = r$; $y_{k_1-1} = y_{k_2+1} = r - 1$. Given such a path y, consider the associated (partially reflected) path y* obtained from y by replacing y_k by $y_k^* = 2r - y_k$ for all $k_1 < k < k_2$, (leaving the other coordinates y_k unchanged). Thus $(y^*)^* = y$, while $y^* = y$ if and only if $k_1 = k_2$.

For each fixed index k with $k_1 \leq k < k_2$, a possible contribution ± 1 to the value $(m_r - n_r)(y^*)$ (for the reflected path y^*), due to a pair $y_k = r$, $y_{k+1} = r \pm 1$, is exactly opposite in sign to the corresponding contribution to the value $(m_r - n_r)(y)$ (for the original path y). Hence, since y and y^* have the same probability 2^{-2n} , one may as well ignore all such contributions, in which case there only remains the single contribution +1 to $(m_r - n_r)(y)$ due to the single pair $y_{k_2} = r$; $y_{k_2+1} = r - 1$. This completes the proof of (3.10).

It now follows from (3.9), (3.10) and (1.4) that

$$C_n(r) = 4D_n(r) = 4P(Y_{2n} = 2r \mid Y_0 = 0) = 4\binom{2n}{n-r}2^{-2n} = 2a_{nr}.$$

Here, we also used the standard André reflection principle. Namely, associate to each path y as above, of length 2n which begins and ends at 0 and meets state r at least once, the path y^* having $y_k^* = 2r - y_k$ when $k \ge k_1$ while $y_k^* = y_k$, otherwise. This sets up a 1:1 correspondence with the set of paths y^* of length 2n which begin at 0 and end at 2r. This completes the proof of Lemma 1.4.

Second proof of Lemma 1.4. Skibinsky (1968); (1969) showed that the mapping from the canonical moments p_i to the power moments c_i is given by the following formulae. Here $q_i = 1 - p_i (i \ge 1)$, $\zeta_i = p_i q_{i-1} (i \ge 1)$ thus $\zeta_1 = p_1$. Define $S_{ij} = 0$ unless $0 \le i \le j$. Further $S_{ij} (0 \le i \le j)$ is recursively defined by $S_{0j} \equiv 1 (j \ge 0)$ and

$$S_{ij} = S_{i,j-1} + \zeta_{j-i+1} S_{i-1,j} \text{ if } 1 \le i \le j.$$
(3.11)

Thus the case j=i reduces to $S_{ii}=\zeta_1S_{i-1,i}$. The moments c_n themselves are finally given by $c_n=S_{nn} (n\geq 0)$. Note that S_{ij} is independent of the p_r with r>j.

For j and n as integers and $n \geq 0$, define

$$Q_j^n = 2^{-n} \binom{n}{m} \text{ if } n = |j| + 2m \text{ with } m = 0, 1, 2, \dots,$$
 (3.12)

and $Q_j^n = 0$ in all other cases. Note from (1.4) that $a_{nr} = 2Q_{2r}^{2n}$. As is easily seen,

$$Q_j^n = \frac{1}{2}(Q_{j-1}^{n-1} + Q_{j+1}^{n-1}) \quad \text{and } Q_{-j}^n = Q_j^n \text{ thus } Q_0^n = Q_1^{n-1}.$$
 (3.13)

Let further S_{ij}^0 denote the value S_{ij} in the special case that $p_k = \frac{1}{2}$ for all $k \geq 1$. Using (3.13), it follows from (3.11) by induction that

$$S_{ij}^{0} = 2^{j-i} Q_{j-i}^{i+j} \quad \text{if } 0 \le i \le j.$$
 (3.14)

For instance $S_{ii}^0 = Q_0^{2i} = Q_1^{2i-1} = \zeta_1 S_{i-1;i}^0$ with $\zeta_1 = p_1 = 1/2$.

Let $r \geq 1$ be fixed, and introduce

$$U_{ij} = 2^{i-j-1} \frac{\partial}{\partial p_r} S_{ij} \mid p_k = 1/2 \text{ for } k \ge 1.$$

Thus $U_{ij} = 0$ unless $0 \le i \le j$ and $r \le j$. Moreover, $U_{0j} \equiv 0$ since $S_{0j} \equiv 1$. We want to show that $\left[\frac{\partial}{\partial p_r}c_n\right]_0 = 2a_{nr}$. In view of $c_n = S_{nn}$ and $a_{nr} = 2Q_{2r}^{2n}$, this is equivalent to $U_{nn} = 2Q_{2r}^{2n}$. More generally, we will show that, for all $0 \le i \le j$.

$$U_{ij} = Q_{j-i+2r}^{i+j} \qquad \text{if } j-i \ge r \ge 1;$$

$$= Q_{j-i+2r}^{i+j} + Q_{i-j+2r}^{i+j} \text{ if } 0 \le j-i < r.$$
(3.15)

For instance $U_{ii} = 2Q_{2r}^{2i}$ and $U_{i-1,i} = Q_{2r+1}^{2i-1} + Q_{2r-1}^{2i-1} (r \ge 2)$; $U_{i-1,i} = Q_3^{2i-1}$ if r = 1.

Differentiating the recursion formula (3.11) with respect to p_r at $p_k = 1/2$ (all $k \ge 1$) and using (3.14), one finds that the U_{ij} satisfy the recursion relation

$$U_{ij} - \frac{1}{2}(U_{i,j-1} + U_{i-1,j}) = -\frac{1}{2}Q_{r+1}^{i+j-1} \quad \text{if } j - i = r;$$

$$= +\frac{1}{2}Q_r^{i+j-1} \quad \text{if } j - i = r - 1;$$

$$= 0 \quad \text{otherwise,}$$

$$(3.16a)$$

as long as $1 \le i < j$. The case j = i is of the form

$$U_{i,i} = U_{i-1,i} + \delta_r^1 Q_1^{2i-1}. \tag{3.16b}$$

The recursion (3.16) and boundary condition $U_{0j} \equiv 0$ together completely determine the U_{ij} . Using (3.13), one easily verifies that $U_{ij} (0 \le i \le j)$ as defined by the right hand side of (3.15), does indeed satisfy (3.16) and $U_{0j} \equiv 0$. This establishes (3.15) and completes the second proof of Lemma 1.4.

Remarks. Formula (3.11) for the S_{ij} , which furnishes a recursive calculation of $c_n = S_{nn}$ from the canonical coordinates p_i , also follows from a simple random walk argument. In fact, the S_{ij} have the simple probabilistic interpretation (3.18) below.

Namely, let $\{X_n\}$ be the random walk on Z_+ described by (3.1), with the p_j as the usual canonical coordinates. We know that $c_n = P_{00}^{(2n)}$, for all $n \geq 0$. Clearly, $P_{0j}^{(n)} = P(X_n = j \mid X_0 = 0)$ satisfy $P_{0j}^{(0)} = \delta_j^0$ and

$$P_{0k}^{(n)} = P_{0;k-1}^{(n-1)} q_{k-1} + P_{0;k+1}^{(n-1)} p_{k+1}, \tag{3.17}$$

 $(n \ge 1; k \ge 0; q_{-1} = 0)$. This allows us to calculate the $P_{0k}^{(n)}$ in a recursive manner. For instance, $c_n = P_{00}^{(2n)} = p_1 P_{01}^{(2n-1)}$. Since $P_{0k}^{(n)} = 0$ if n < k, (3.17) is trivially satisfied when

n < k. Also note that $P_{0k}^{(k)} = q_0 q_1 \dots q_{k-1}$. All terms in (3.17) vanish unless n = k + 2i with $i \in \mathbb{Z}_+$, in which case n = i + j; k = j - i with $0 \le i \le j$ as integers. It follows from (3.17) that the S_{ij} defined by

$$S_{ij} = \frac{1}{q_0 q_1 q_2 \dots q_{j-i-1}} P_{0;j-i}^{(i+j)} \text{ for } 0 \le i \le j,$$
(3.18)

 $(q_0=1)$ satisfy the recursion relation (3.11). Moreover, $S_{0k}=P_{0k}^{(k)}/q_0q_1\ldots q_{k-1}=1$, for all $k\geq 0$. Finally, $c_n=P_{00}^{(2n)}=S_{nn}$.

In view of the interpretation (3.18) of the S_{ij} , formula (3.15) can also be regarded as an explicit formula for the quantities $\left[\frac{\partial}{\partial p_r}P_{0j}^{(n)}\right]_0$, equivalently, as an explicit formula for $E[(m_r - n_r)(X_n = j) \mid X_0 = 0]$, with m_r , n_r as in (3.7).

Theorem 2 in Skibinsky (1968) also has a simple probabilistic proof. It states that

$$c_n = \sum_{0 \le i \le n/2} (S_{i,n-i})^2 \prod_{j=1}^{n-2i} \zeta_j.$$
 (3.19)

In fact, paying attention to the value $X_n = k$ (say),

$$c_n = P(X_{2n} = 0 \mid X_0 = 0) = \sum_k P_{0k}^{(n)} P_{k0}^{(n)} = \sum_k \frac{1}{\pi_k} (P_{0k}^{(n)})^2.$$
 (3.20)

Here, $\pi_k = q_0 q_1 \dots q_{k-1}/p_1 p_2 \dots p_k$, $(\pi_0 = 1)$. We also used the well known relation $\pi_i P_{ij}^{(n)} = \pi_j P_{ji}^{(n)}$, (all i, j, n; see for instance Karlin and McGregor (1959) p. 68). Noting that $P_{0k}^{(n)}$ vanishes unless k = n - 2i with $0 \le i \le n/2$, and using (3.18), one easily verifies that (3.19), (3.20) are equivalent.

4. Further proofs

Proof of Lemma 2.1. Let $i, j \ge 1$. From $\Sigma = \frac{1}{2}AA'$ and $a_{kr} = 0$ for r > k, one has

$$\sigma_{ij} = \frac{1}{2} \sum_{r=1}^{\min(i,j)} a_{ir} a_{jr} = 2^{-2i-2j+1} \sum_{r=1}^{\min(i,j)} {2i \choose i-r} {2j \choose j-r}$$

$$= -c_i^0 c_j^0 + \sum_{r=-\min(i,j)}^{\min(i,j)} 2^{-2i} {2i \choose i-r} 2^{-2j} {2j \choose j+r} = -c_i^0 c_j^0 + c_{i+j}^0,$$

proving (2.1). After all, the latter sum is equal to the coefficient of z^{i+j} in the expansion of $\left(\frac{1+z}{2}\right)^{2i}\left(\frac{1+z}{2}\right)^{2j}$.

Recall that $c_k^0 \approx 1/\sqrt{\pi k}$ as $k \to \infty$. Hence, $\sigma_{jj} = c_{2j}^0 - (c_j^0)^2 \approx (2\pi j)^{-1/2}$ and $\sigma_{ij} = c_{i+j}^0 - c_i^0 c_j^0 \approx (1 - c_i^0)(\pi j)^{-1/2}$ as $j \to \infty$. Thus, for i fixed and $j \to \infty$,

$$\rho_{ij} = \frac{\sigma_{ij}}{\sqrt{\sigma_{ii}}\sqrt{\sigma_{jj}}} \approx D_i \ j^{-1/4}, \text{ where } D_i = (\pi/2)^{-1/4} (1 - c_i^0)(\sigma_{ii})^{-1/2}.$$

In particular $\rho_{s,s+r} \to 0$ as $r \to \infty$. If both i and j tend to infinity then

$$\sigma_{ij} = c_{i+j}^0 (1 - c_i^0 c_j^0 / c_{i+j}^0) \approx c_{i+j}^0 \approx 1 / \sqrt{\pi(i+j)}.$$

Here we used that $c_i^0 c_j^0/c_{i+j}^0 \approx \left[\frac{1}{\pi} \left(\frac{1}{i} + \frac{1}{j}\right)\right]^{1/2} \to 0$. Hence, if $i, j \to \infty$ in such a way that $j/i \to K$ then

$$\rho_{ij} \approx \left[\frac{4ij}{(i+j)^2}\right]^{1/4} \to \left(\frac{4K}{(K+1)^2}\right)^{1/4}.$$

Proof of (2.7). We want to prove that the coefficients b_{im} in (2.5) are given by (2.7). Letting $y = \cos \theta = 2x - 1$, one has $\cos n\theta = T_n(y) = T_n^*(x)$ thus

$$\sum_{n=0}^{\infty} T_n^*(x) u^n = \sum_{n=0}^{\infty} \cos n\theta u^n = Re \left[\sum_{n=0}^{\infty} (e^{i\theta} u)^n \right] = Re \frac{1}{1 - ue^{i\theta}}$$
$$= \frac{1 - u \cos \theta}{1 + u^2 - 2u \cos \theta} = \frac{1 + u - 2ux}{(1 + u)^2 - 4ux} = (1 + u - 2ux) \sum_{r=0}^{\infty} (4ux)^r (1 + u)^{-2r-2}.$$

The coefficient of x^m is found to be $2^{2m-1}u^m(1-u)(1+u)^{-2m-1}$. Expanding the latter in powers of u, we find that the coefficient of u^n is precisely b_{nm} as given by (2.7).

Proof of the identity (2.11). This identity must be known. Recall that $T_r^*(x) = \cos r\theta$ when $x = (\cos \frac{\theta}{2})^2$. If $m \ge 1$ then

$$x^{m} = \left(\cos\frac{\theta}{2}\right)^{2m} = 2^{-2m} \left(e^{i\theta/2} + e^{-i\theta/2}\right)^{2m} = 2^{-2m} \sum_{j=0}^{2m} {2m \choose j} \cos(m-j)\theta.$$

The term with j=m gives rise to $2^{-2m}\binom{2m}{m}=c_m^0$. Further, for $r=1,\ldots,m$, the two terms with $j=m\pm r$ together give rise to $2^{-2m+1}\binom{2m}{m-r}\cos r\theta=a_{mr}T_r^*(x)$, in view of (1.4). This proves (2.11).

Proof that AB = I, see Theorem 2.3. Here A, B are lower triangular hence also C = AB. Further $c_{ii} = a_{ii}b_{ii} = 1$ thus it suffices to show that $c_{im} = 0$ when $1 \le m < i$. From (1.4) and (2.7),

$$c_{im} = \sum_{j=m}^{i} a_{ij} b_{jm} = \sum_{j=m}^{i} 2^{-2i+1} {2i \choose i+j} (-1)^{j+m} 2^{2m-1} \frac{j}{m} {j+m-1 \choose 2m-1}.$$

This can be written as $c_{im} = \sum_{j=m}^{i} (-1)^{j} {2i \choose i+j} g(j)$, where

$$g(x) = \alpha x \binom{x+m-1}{2m-1} = \frac{\alpha x^2}{(2m-1)!} \prod_{r=1}^{m-1} (x+r)(x-r),$$

with $\alpha = \alpha_{im}$ as a constant factor. Note that g(x) is an even polynomial of degree 2m such that g(r) = 0 for $r = 0, \pm 1, \ldots, \pm (m-1)$. Hence, letting i + j = s,

$$2c_{im} = \sum_{j=-i}^{i} (-1)^{j} {2i \choose i+j} g(j) = \sum_{s=0}^{2i} (-1)^{s-i} {2i \choose s} g(s-i) = (-1)^{i} \Delta^{2i} g(-i) = 0,$$

since g is of degree 2m < 2i. Here $\Delta = E - 1$ is the usual difference operator thus (Eg)(x) = g(x+1).

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