INFINITE PRODUCT EXPANSIONS FOR MATRIX n-th ROOTS

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1. Introduction

In this paper a denotes a square matrix with real or complex elements (though the theorems and their proofs are valid in any Banach algebra). Its spectral radius $\rho(a)$ is given by

(1)
$$\rho(a) = \lim ||a^{\nu}||^{1/\nu}, \text{ as } \nu \to \infty,$$

with any matrix norm (see [4], p. 183). If $\rho(a) < 1$ and n is a positive integer then the binomial series

(2)
$$S(a) = \sum_{\nu=0}^{\infty} {\binom{-1/n}{\nu}} (-a)^{\nu}$$

converges and its sum satisfies $S(a)^n = (1-a)^{-1}$. Let

(3)
$$u(x) = 1 + \sum_{\nu=1}^{q-1} \frac{\Gamma(n^{-1} + \nu)x^{\nu}}{\nu ! \Gamma(n^{-1})},$$

where q is any integer exceeding 1. Then u(a) is the sum of the first q terms of the series (2). Write

(4)
$$f(x) = 1 + u(x)^{n}(x-1)$$

and let a_0 , a_1 , a_2 , \cdots be the sequence of matrices obtained by the iterative procedure

(5)
$$a_0 = a, \quad a_{\nu+1} = f(a_{\nu}).$$

Defining polynomials $\phi_0(x)$, $\phi_1(x)$, $\phi_2(x)$, \cdots inductively by

(6)
$$\phi_0(x) = x, \quad \phi_{\nu+1}(x) = f(\phi_{\nu}(x)),$$

we have $a_{\nu} = \phi_{\nu}(a)$ and therefore $a_{\mu}a_{\nu} = a_{\nu}a_{\mu}$ for all μ , ν . The following is proved in section 2:

Theorem 1. If $\rho(a) < 1$ then

(7)
$$P(a) = \prod_{\nu=0}^{\infty} u(a_{\nu})$$

converges and P(a) = S(a). Furthermore, if $\rho(a) < r < 1$, then

$$||a_{\nu}|| < Mr^{q^{\nu}}$$

for all v, where M depends on r and a but is independent of v and q.

Inequality (8) shows that P(a) converges very rapidly. This could make it useful for the numerical computation of S(a). In general the series (2) converges too slowly to be used for this purpose. In section 3 it is shown that when n > 1 the infinite product (7) converges for a larger class of matrices a than does the series (2). If n = 1 then (3) and (4) give $f(x) = x^{a}$. The solution of (5) is then $a_{n} = a^{a^{n}}$ and (7) reduces to

(9)
$$(1-a)^{-1} = \prod_{\nu=0}^{\infty} \{1 + a^{q\nu} + a^{2q\nu} + \cdots + a^{(q-1)q^{\nu}}\}.$$

This well-known formula goes back to Euler. Its use for practical computation was suggested by Ostrowski [6], Hotelling [3] and others. Hotelling was able to connect (9) in the special case q=2 with an iterative method for matrix inversion given by the Newton-Raphson formula. For (7) there is a similar connection with the Newton-Raphson formula which is discussed in section 4. Theorem 1 can be used to find a matrix c satisfying $c^n = b$ for any square matrix b whose spectrum lies entirely in the half plane $\text{Re } \lambda > 0$. For the spectrum of $a = (b+1)^{-1}(b-1)$ then lies in the disc $|\lambda| < 1$ and $c = (1+a)^{1/n}(1-a)^{-1/n}$ can be computed with the help of (7). In the special case when the eigenvalues of b are real and positive it is simpler to take $a = 1 - k^{-1}b$, where b is any real number satisfying $b > \frac{1}{2}\rho(b)$. The eigenvalues of b then satisfy $-1 < \lambda < 1$ and $c = k^{1/n}(1-a)^{1/n}$ can be computed with the help of (7).

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LEMMA 1. If m is any integer in the range $1 \le m \le n$ and

$$u(x)^m = \beta_0 + \beta_1 x + \beta_2 x^2 + \cdots$$

then for all v,

$$\beta_{\nu} \geq \beta_{\nu+1} \geq 0.$$

PROOF. Since $(1-x)^{-1/n} = u(x) + x^q v_1(x)$, it follows that

(11)
$$(1-x)^{-m/n} = u(x)^m + x^q v_m(x),$$

where $v_1(x)$, $v_m(x)$ are power series with positive coefficients. Comparing coefficients in (11) we get $(-1)^{\nu} {-m/n \choose \nu} = \beta_{\nu}$ for $0 \le \nu < q$. If $1 \le m \le n$ then $(-1)^{\nu} {-m/n \choose \nu}$ is a positive monotonic decreasing function of ν . Hence, (10) holds in the range $0 \le \nu < q-1$ and its remains to prove (10) for the

range $\nu \ge q-1$. We do this by induction over m. Clearly (10) holds when m=1 because then $\beta_{\nu}=0$ for all $\nu \ge q$. If the lemma holds for some m in the range $1 \le m < n$ then (3) gives

$$u(x)^{m+1} = u(x)^m u(x) = \left(\sum_{\nu=0}^{\infty} \beta_{\nu} x^{\nu}\right) \left(\sum_{\nu=0}^{q-1} \alpha_{\nu} x^{\nu}\right) = \sum_{\nu=0}^{\infty} \gamma_{\nu} x^{\nu},$$

where

$$\alpha_{\mu} = \frac{\Gamma(n^{-1} + \mu)}{\mu ! \Gamma(n^{-1})}$$
 and $\gamma_{\nu} = \sum_{\mu=0}^{q-1} \alpha_{\mu} \beta_{\nu-\mu}$

for all $\nu \ge q-1$. Since (10) holds for all ν we have

$$0 \leq \gamma_{\nu+1} = \sum_{\mu=0}^{q-1} \alpha_{\mu} \beta_{\nu+1-\mu} \leq \sum_{\mu=0}^{q-1} \alpha_{\mu} \beta_{\nu-\mu} = \gamma_{\nu},$$

for all $v \ge q-1$. The lemma is therefore true for m+1 also. This establishes Lemma 1.

LEMMA 2. $f(x) = x^q g(x)$ and $\phi_{\nu}(x) = x^{q\nu} \psi_{\nu}(x)$ where g(x), $\psi_{\nu}(x)$ are polynomials with real non-negative coefficients which satisfy $g(1) = \psi_{\nu}(1) = 1$.

Proof. With m = n, (11) gives

$$x^q v_n(x)(1-x) = 1 + u(x)^n(x-1) = f(x).$$

Hence $f(x) = x^q g(x)$ for some polynomial g(x). If $u(x)^n = \beta_0 + \beta_1 x + \beta_2 x^2 + \cdots$ then $\beta_0 = 1$ and $x^q g(x) = 1 + u(x)^n (x-1) = \sum_{\nu=0}^{\infty} (\beta_{\nu} - \beta_{\nu+1}) x^{\nu}$.

The coefficients of g(x) are therefore non-negative by Lemma 1. Also g(1) = f(1) = 1 from (4). Induction over ν will be used to prove that $\phi_{\nu}(x)$ is of the form $x^{q\nu} \psi_{\nu}(x)$, where $\psi_{\nu}(x)$ has non-negative coefficients. This is trivial for $\nu = 0$ since $\phi_0(x) = x$. If it is true for some integer ν then (6) gives

$$\phi_{\nu+1} = (\phi_{\nu})^q g(\phi_{\nu}) = (x^{q\nu}\psi_{\nu})^q g(\phi_{\nu}) = x^{q^{\nu+1}}\psi_{\nu+1},$$

where $\psi_{\nu+1} = (\psi_{\nu})^q g(\phi_{\nu})$ is a polynomial with non-negative coefficients. The result is therefore true for $\nu+1$ also and the induction is complete. Since f(1)=1 it follows from (6) by induction that $\phi_{\nu}(1)=1$ for all ν . Hence $\psi_{\nu}(1)=\phi_{\nu}(1)=1$ and the proof of Lemma 2 is finished.

When q = 2, (3) gives $u(x) = 1 + n^{-1}x$. Then (4) gives

(12)
$$f = 1 + nu^{n+1} - (n+1)u^{n},$$

$$= (u-1)^{2}(1 + 2u + 3u^{2} + \cdots + nu^{n-1}),$$

$$= x^{2}n^{-2}(1 + 2u + 3u^{2} + \cdots + nu^{n-1}).$$

This is the relation $f(x) = x^q g(x)$ for the special case q = 2. If x is small then f(x) can be computed more accurately from the relation $f(x) = x^q g(x)$ than from (4) which involves internal cancellation when x is small.

Proof of theorem 1. If $\rho(a) < r$ then (1) gives

$$(13) ||a^{\nu}|| \leq Mr^{\nu}$$

for some constant M. Hence $||\phi_{\nu}(a)|| \leq M\phi_{\nu}(r)$ since the coefficients of $\phi_{\nu}(x)$ are non-negative by Lemma 2. Since 0 < r < 1, Lemma 2 also gives $||\phi_{\nu}(a)|| \leq Mr^{q\nu}\psi_{\nu}(r) \leq Mr^{q\nu}\psi_{\nu}(1) = Mr^{q\nu}$.

This proves (8) because $a_{\nu} = \phi_{\nu}(a)$. From (4) and (5),

$$(1-a_{\nu+1})=u(a_{\nu})^n(1-a_{\nu}).$$

Since $a_{\mu}a_{\nu}=a_{\nu}a_{\mu}$ it follows by induction that

(14)
$$(1-a_{\kappa+1}) = \left[\prod_{\nu=0}^{\kappa} u(a_{\nu})\right]^{n} (1-a).$$

Since $a_{\nu} = \phi_{\nu}(a)$ we have $\prod_{\nu=0}^{\kappa} u(a_{\nu}) = W_{\kappa}(a)$, where $W_{\kappa}(x)$ is a polynomial with non-negative coefficients. When -1 < x < 1 it follows from (2) and (14) that

(15)
$$S(x) = (1-x)^{-1/n} = W_{\kappa}(x)\{1-\phi_{\kappa+1}(x)\}^{-1/n}.$$

Expanding $\{1-\phi_{\kappa+1}(x)\}^{-1/n}$ by the binomial theorem and using Lemma 2 we get

$$S(x) = W_{\kappa}(x) + x^{q^{\kappa+1}} V_{\kappa}(x),$$

where $V_{\kappa}(x)$ is a power series with non-negative coefficients. This and (13) give $||S(a)-W_{\kappa}(a)||=||a^{q^{\kappa+1}}V_{\kappa}(a)|| \leq Mr^{q^{\kappa+1}}V_{\kappa}(r)$. Since none of the coefficients of S(x) exceeds 1 by (2), the same is true of the coefficients of $V_{\kappa}(x)$ by (16). Hence,

(17)
$$||S(a) - W_{\kappa}(a)|| \leq M r^{q^{\kappa+1}} (1-r)^{-1},$$

$$S(a) = \lim_{r \to \infty} W_{\kappa}(a) = \prod_{\nu=0}^{\infty} u(a_{\nu}).$$

This completes the proof of Theorem 1.

3. Domain of convergence

Let $D = \bigcup_{\nu=0}^{\infty} D_{\nu}$ where D_{ν} is the set of points in the complex z plane for which $|\phi_{\nu}(z)| < 1$. Each D_{ν} is an open set and D_{0} is the disc |z| < 1.

THEOREM 2. If the spectrum of a lies wholly in D then the infinite product (7) converges and satisfies $P(a)^n = (1-a)^{-1}$.

PROOF. Lemma 2 gives |f(z)| < g(|z|) < g(1) = 1 for |z| < 1. This and (6) show that $|\phi_{\nu+1}(z)| < 1$ when $|\phi_{\nu}(z)| < 1$. That is,

 $D_{\nu+1} \supset D_{\nu}$ for all ν . Since the spectrum of a is compact it must lie in D_{μ} for some μ . Then the spectrum of $a_{\mu} = \phi_{\mu}(a)$ lies wholly in the disc |z| < 1. Hence $\rho(a_{\mu}) < 1$ and $S(a_{\mu}) = \prod_{\nu=\mu}^{\infty} u(a_{\nu})$ by Theorem 1. From (14) with $\kappa = \mu - 1$ we get

$$P(a)^n = S(a_{\mu})^n \left[\prod_{\nu=0}^{\mu-1} u(a_{\nu}) \right]^n = S(a_{\mu})^n (1-a_{\mu}) (1-a)^{-1} = (1-a)^{-1}.$$

This completes the proof of Theorem 2. If n=1 then $\phi_{\nu}(z)=z^{q^{\nu}}$ and $D_{\nu}=D_0$ for all ν . That this is not so when n>1 is shown by the next theorem.

THEOREM 3. If n > 1 then D_1 includes all of the closed disc $|z| \leq 1$ except the point z = 1.

PROOF. Lemma 2 gives $|f(z)| \le f(|z|) \le f(1) = 1$ for $|z| \le 1$. The inequality $|f(z)| \le f(|z|)$ can reduce to equality only when the terms of f(z) all have the same complex argument (see [2], p. 26). If

$$u(z) = k_1 z^{q-1} + k_2 z^{q-2} + \cdots$$

then (4) gives

$$f(z) = k_1^n z^{1+n(q-1)} + (nk_2 - k_1)k_1^{n-1} z^{n(q-1)} + \cdots,$$

where the terms shown are those of the highest degrees. If n>1 then both these terms have positive coefficients because $k_2 \ge k_1 > 0$ by (3). These terms have the same complex argument only when z is real and positive. Therefore $|f(z)| \le f(|z|)$ reduces to equality only when z is real and positive. Hence z=1 is the only point of the disc $|z| \le 1$ at which $|f(z)| \le 1$ reduces to equality. Since $f(z) = \phi_1(z)$ it follows that D_1 includes all of the disc $|z| \le 1$ except the point z=1. This established Theorem 3.

Since D_1 is an open set Theorem 3 shows that a part of it must lie outside the circle |z|=1 when n>1. The region of convergence of the infinite product P(a) is therefore somewhat larger than that of the series (2) which diverges if a has an eigenvalue outside the circle |z|=1. More precise information about the size of D will be given only for the special case q=2. Let H be the convex hull of the set which is the union of the closed disc $|z| \leq 1$ and the single point z=-n. Let H_0 be the set obtained by deleting from H the two points z=-n, 1.

Theorem 4. If q=2 and n>1 then D_1 includes H_0 .

PROOF. If $z+n = de^{i\delta}$ then 0 < d < n+1 and

$$-\frac{1}{n} \le \sin \delta \le \frac{1}{n}$$

for all z in H_0 . Also $u(z) = 1 + n^{-1}z = n^{-1}de^{i\delta}$ since q = 2. From (12) we get

$$|f(z)|^2 = \{1 + nu^{n+1} - (n+1)u^n\}\{1 + n\bar{u}^{n+1} - (n+1)\bar{u}^n\}.$$

With $u = n^{-1} de^{i\delta}$ this gives $|f(z)|^2 = 1 + (d/n)^n h(d, \delta)$ where

$$h(d, \delta) = (d/n)^n \{d^2 - 2(n+1)d\cos\delta + (n+1)^2\} + 2d\cos(n+1)\delta - 2(n+1)\cos n\delta.$$

Hence |f(z)| < 1 if and only if $h(d, \delta) < 0$. Since $f(z) = \phi_1(z)$ it follows that $z \in D_1$ if and only if $h(d, \delta) < 0$. To prove Theorem 4 it is therefore sufficient to show that $h(d, \delta) < 0$ throughout H_0 . Clearly $h(d, \delta) < 0$ in $H_0 \cap \bar{D}_0$ since D_1 includes the closed disc \bar{D}_0 , except for the point z = 1, by Theorem 3. To prove that $h(d, \delta) < 0$ in the whole of H_0 it is therefore sufficient to show that $\partial h/\partial d > 0$ in $H_0 - \bar{D}_0$ because each point of $H_0 - \bar{D}_0$ lies on some line segment joining z = -n to a point of $H_0 \cap \bar{D}_0$. Since $|z|^2 = |-n + de^{i\delta}|^2 = n^2 - 2nd \cos \delta + d^2$, we can express $\partial h/\partial d$ in the form

Since $x^{-1} \sin x \ge 3/\pi$ in $0 < x \le \pi/6$, (18) gives

$$\sin |\delta| \le n^{-1} \le 3(2n+2)^{-1} \le \sin \{(2n+2)^{-1}\pi\}.$$

Hence $|\delta| \le (2n+2)^{-1}\pi$ and $\cos (n+1)\delta \ge 0$ for all z in H_0 . This and (19) give

$$\partial h/\partial d \ge n^{-(n+1)} d^{n-1} \{ (n+1)^2 |z|^2 - d^2 \} > n^{-(n+1)} d^{n-1} (n+1)^2 \{ |z|^2 - 1 \},$$

for all z in H_0 . Therefore $\partial h/\partial d > 0$ in $H_0 - \bar{D}_0$. This completes the proof of Theorem 4. Notice that the points z = -n, 1 which were omitted from H_0 lie outside D because f(-n) = 1 when q = 2 and therefore $\phi_{\nu}(-n) = \phi_{\nu}(1) = 1$ for all $\nu \ge 1$ by (6).

4. Connection with Newton-Raphson

The Newton-Raphson formula for the numerical solution of an equation Y(x) = 0 is $x_{\nu} - x_{\nu+1} = Y(x_{\nu})/Y'(x_{\nu})$. With $Y(x) = b - x^{-n}$ this becomes

(20)
$$x_{\nu+1} = x_{\nu} \{ 1 + n^{-1} (1 - b x_{\nu}^{n}) \}.$$

As a generalisation of this we consider the formula

(21)
$$x_{\nu+1} = x_{\nu} u(a_{\nu}), \quad a_{\nu} = 1 - b x_{\nu}^{n},$$

where u(x) is given by (3) with any q. This reduces to (20) when q=2 because then $u(x)=1+n^{-1}x$. When x_{ν} and b are square matrices, (21) gives

(22)
$$x_{\nu+1} = x_0 u(a_0) u(a_1) \cdots u(a_{\nu}).$$

Also, (4) and (21) show that $f(a_{\nu}) = 1 + (a_{\nu} - 1)u(a_{\nu})^n = 1 - bx_{\nu}^n u(a_{\nu})^n$.

If $x_{\nu}a_{\nu} = a_{\nu}x_{\nu}$ then $x_{\nu}^{n}u(a_{\nu})^{n} = x_{\nu+1}^{n}$ by (21) and

(23)
$$f(a_{\nu}) = 1 - b(x_{\nu+1})^n = a_{\nu+1}.$$

Compare this with (5). The condition $x_{\nu}a_{\nu}=a_{\nu}x_{\nu}$ is satisfied if $x_{\nu}b=bx_{\nu}$ and this is true by induction provided that $x_{0}b=bx_{0}$. When this is so, (22), (23) and Theorem 1 show that $x_{\nu+1}\to x_{0}P(a_{0})$ as $v\to\infty$ provided that $\rho(a_{0})<1$. If $x_{0}P(a_{0})=L$ then

$$L^n = x_0^n P(a_0)^n = x_0^n (1-a_0)^{-1} = b^{-1}.$$

The following theorem is therefore true.

THEOREM 5. If $a_0 = 1 - bx_0^n$ has $\rho(a_0) < 1$ and $x_0b = bx_0$ then the sequence of matrices x_0, x_1, x_2, \cdots obtained from (21) tends to a limit matrix L which satisfies $L^n = b^{-1}$. Furthermore, the rate of convergence is of the q-th order.

Altman [1] and Petryshyn [7] have studied (21) in the special case when n=1. They obtain results similar to Theorem 5 but without the requirement $x_0b=bx_0$. This requirement can be deleted from Theorem 5 in the case n=1 because (23) then follows without use of the relation $x_{\nu}a_{\nu}=a_{\nu}x_{\nu}$. The following counter-example shows that $x_0b=bx_0$ cannot be deleted from Theorem 5 when n>1. If

$$b=egin{pmatrix} 1 & 0 \ 0 & \mu^{-n} \end{pmatrix}$$
 , $x_
u=egin{pmatrix} 1 & \xi_
u \ 0 & \mu \end{pmatrix}$, $x_
u^n=egin{pmatrix} 1 & \zeta\xi_
u \ 0 & \mu^n \end{pmatrix}$

then $\zeta = (\mu - 1)^{-1}(\mu^n - 1)$ and $\rho(a_0) = 0$ where $a_0 = 1 - bx_0^n$. These matrices satisfy (20) provided that $\xi_{\nu} = (1 - n^{-1}\zeta)^{\nu}\xi_0$ for all ν . If n > 1 and $\mu \ge 4$ then $|1 - n^{-1}\zeta| > 1$ and x_{ν} does not tend to a limit as $\nu \to \infty$ because $|\xi_{\nu}| \to \infty$. The deletion of $x_0 b = bx_0$ from Theorem 5 therefore produces a false proposition when n > 1. When n > 1 and q = 2, Theorem 4 enables the condition $\rho(a_0) < 1$ in Theorem 5 to be replaced by the requirement that the spectrum of a_0 lie wholly in H_0 .

With $Y(x) = bx^n - 1$ the Newton-Raphson formula becomes

(24)
$$x_{\nu+1} = (1-n^{-1})x_{\nu} + (nbx_{\nu}^{n-1})^{-1}.$$

A higher order formula of Traub [8] generalises this is the same way that (21) generalises (20). Laasonen [5] has shown that (24) can be used to find a square root of any real matrix whose eigenvalues are all positive. Each iteration of (24) requires a matrix inversion which could introduce considerable error when $b^{-1/n}$ is ill-conditioned. The iterations of (5) and (21) do not involve matrix inversions.

References

- [1] M. Altman, 'An optimum cubically convergent iterative method for inverting a linear bounded operator in Hilbert space', Pacific J. Math. 10 (1960), 1107-1113.
- [2] G. H. Hardy, J. E. Littlewood and G. Polya, Inequalities (C.U.P., 2nd ed. 1952).
- [3] H. Hotelling, 'Some new methods in matrix calculation', Ann. Math. Statist. 14 (1943), 1-34.
- [4] A. S. Householder, The theory of matrices in numerical analysis (Blaisdell, New York, 1964).
- [5] P. Lassonen, 'On the iterative solution of the matrix equation $AX^2-I=O$ ', Math. Tables Aids Comput. 12 (1958), 109-116.
- [6] A. Ostrowski, 'Sur quelques transformations de la série de Liouville-Neumann', C.R. Acad. Sci. Paris 206 (1938), 1345-1347.
- [7] W. V. Petryshyn, 'On the inversion of matrices and linear operators', Proc. Amer. Math. Soc. 16 (1965), 893-901.
- [8] J. F. Traub, 'Comparison of iterative methods for the calculation of n-th roots', Comm. ACM 4 (1961), 143-145.

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