Sensitivity of finite Markov chains under perturbation

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Abstract: Meyer (1992) has developed inequalities in terms of the non-unit eigenvalues λ_i , $j=2,\ldots,n$, of a stochastic matrix P containing a single irreducible set of states, for the condition number $\max |a_{ij}^{\#}|$, where $A^{\#} = \{a_{ij}^{\#}\}$ is the group generalized inverse of A = I - P. In this note we derive, succinctly, analogous inequalities for the alternative condition number, the ergodicity coefficient $\tau_1(A^{\#})$, using the properties of ergodicity coefficients: $(\min |1 - \lambda_i|)^{-1} \le \tau_1(A^{\#}) \le \sum (1 - \lambda_i)^{-1}$.

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1. Introduction

If P is an $n \times n$ stochastic matrix containing a single irreducible set of states, so that there is a unique stationary distribution vector $\boldsymbol{\pi}^T = \{\pi_i\}(\boldsymbol{\pi}^T P = \boldsymbol{\pi}^T, \, \boldsymbol{\pi}^T \mathbf{1} = 1)$, and \bar{P} is any other $n \times n$ stochastic matrix, with $\bar{\boldsymbol{\pi}}^T = \{\bar{\pi}_i\}$ a stationary distribution vector corresponding to it, then direct matrix manipulation yields

$$(\overline{\boldsymbol{\pi}}^{\mathrm{T}} - \boldsymbol{\pi}^{\mathrm{T}})(I - P + 1\boldsymbol{\pi}^{\mathrm{T}}) = \overline{\boldsymbol{\pi}}^{\mathrm{T}}E$$

where $E=\{e_{ij}\}=\overline{P}-P$. Then, denoting by $Z=\{z_{ij}\}=(I-P+\mathbf{1}\boldsymbol{\pi}^{\mathrm{T}})^{-1}$ the fundamental matrix (Kemeny and Snell, 1960) of the Markov chain governed by P, and by $A^{\#}=\{a_{ij}^{\#}\}=(I-P+\mathbf{1}\boldsymbol{\pi}^{\mathrm{T}})^{-1}-\mathbf{1}\boldsymbol{\pi}^{\mathrm{T}}=Z-\mathbf{1}\boldsymbol{\pi}^{\mathrm{T}}$ the group generalized inverse of $A=\{a_{ij}\}=I-P$, we obtain from (1), since $\overline{\boldsymbol{\pi}}^{\mathrm{T}}E\mathbf{1}=0$,

$$\bar{\boldsymbol{\pi}}^{\mathrm{T}} - \boldsymbol{\pi}^{\mathrm{T}} = \bar{\boldsymbol{\pi}}^{\mathrm{T}} E A^{\#}. \tag{2}$$

Thus

$$\left| \overline{\pi}_{j} - \pi_{j} \right| = \left| \sum_{i} \sum_{s} \overline{\pi}_{i} e_{is} a_{sj}^{\#} \right| \leq \left(\sum_{i} \sum_{s} \overline{\pi}_{i} \left| e_{is} \right| \right) \left(\max_{k,j} \left| a_{kj}^{\#} \right| \right) \leq \left(\sum_{i} \overline{\pi}_{i} \right) \left(\max_{r} \sum_{s} \left| e_{rs} \right| \right) \left(\max_{k,j} \left| a_{kj}^{\#} \right| \right)$$

so that (Funderlic and Meyer, 1986)

$$\max_{j} \left| \overline{\pi}_{j} - \pi_{j} \right| \leqslant \|E\|_{1} \max_{i,j} \left| a_{ij}^{\#} \right|. \tag{3}$$

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(Equation (1) and equations close to (2) and (3) first occur in Schweitzer (1968).) On the other hand, from (2),

$$\| \overline{\boldsymbol{\pi}}^{T} - \boldsymbol{\pi}^{T} \|_{1} = \| \overline{\boldsymbol{\pi}}^{T} E \|_{1} (\| \overline{\boldsymbol{\pi}}^{T} E A^{\#} \|_{1} / \| \overline{\boldsymbol{\pi}}^{T} E \|_{1}) \leq \| \overline{\boldsymbol{\pi}}^{T} \|_{1} \| E \|_{1} \sup_{\| \boldsymbol{\delta}^{T} \|_{1} = 1, \boldsymbol{\delta}^{T} \mathbf{1} = 0} \| \boldsymbol{\delta}^{T} A^{\#} \|_{1}$$

so that (Seneta, 1988, 1991)

$$(\|\bar{\boldsymbol{\pi}}^{\mathrm{T}} - \boldsymbol{\pi}^{\mathrm{T}}\|_{1} / \|\boldsymbol{\pi}^{\mathrm{T}}\|_{1}) / (\|E\|_{1} / \|P\|_{1}) \leq \tau_{1}(A^{\#})$$

$$(4)$$

since $\| \overline{\boldsymbol{\pi}}^{T} \|_{1} = \| \boldsymbol{\pi}^{T} \|_{1} = \| P \|_{1} = 1$.

For any $n \times n$ matrix $B = \{B_{ij}\}$ with equal row sums b (i.e. B1 = b1),

$$\tau_1(B) \stackrel{\text{def}}{=} \sup_{\|\boldsymbol{\delta}^{\mathrm{T}}\| = 1, \boldsymbol{\delta}^{\mathrm{T}} \mathbf{1} = 0} \|\boldsymbol{\delta}^{\mathrm{T}} B\|_1 = \max_{i,j} \frac{1}{2} \sum_{s=1}^n |B_{is} - B_{js}|$$
 (5)

(Seneta, 1984) and further, for any eigenvalue λ of B such that $\lambda \neq b$, it is true that

$$|\lambda| \leqslant \tau_1(B) \tag{6}$$

providing there is a left eigenvector v^{T} corresponding to b such that $v^{T}\mathbf{1} = 1$ (Seneta, 1984). This property holds for the matrices B = P, Z, I - P, $A^{\#}$ with $v^{T} = \pi^{T}$ and b = 1, 1, 0, 0 respectively. These properties will be useful in the sequel.

On the basis of (3), Funderlic and Meyer (1986), Meyer (1992) choose

$$\kappa(A) = \max_{i,j} \left| a_{ij}^{\#} \right|$$

as a measure of relative sensitivity ('condition number') of π under perturbation of P to \overline{P} , while on the basis of (4) Seneta (1991) proposes $\tau_1(A^{\#}) \equiv \tau_1(Z)$ as condition number.

Recently, Meyer (1992) has addressed the question whether the closeness of the non-unit eigenvalues of P to unit provides complete information about the relative sensitivity of P. He has established that this is so by deriving the inequalities

$$\frac{1}{n \min_{2 \le i \le n} |1 - \lambda_i|} \le \kappa(A) \le \frac{2(n-1)\delta}{\prod_{j=2}^n (1 - \lambda_j)} \tag{7}$$

where $1, \lambda_2, \ldots, \lambda_n$ are the eigenvalues of P. Here

$$\delta = \max_{i,j,i\neq j} \prod_{k\neq i,j} a_{kk}$$

is the product of all but the two smallest diagonal entries of A = I - P. Note that since Z has the eigenvalues 1, $(1 - \lambda_i)^{-1}$, j = 2, ..., n,

$$\kappa(A) \leqslant 2(n-1)\delta \det Z. \tag{8}$$

We note that since the eigenvalues of $A^{\#}$ are 0, $(1 - \lambda_i)^{-1}$, j = 2, ..., n, and $A^{\#}\mathbf{1} = \mathbf{0}$, by (6),

$$\frac{1}{\min_{2 \leqslant j \leqslant n} |1 - \lambda_j|} \leqslant \tau_1(A^\#) \tag{9}$$

which is a slight sharpening of an inequality in Seneta (1991), produced at Meyer's instigation.

Thinking of $\tau_1(A^{\#})$ as condition number, we note that (9) has a sharper form than the left-hand side of (7). The question arises whether, using the properties (5) and (6) of the 'coefficient of ergodicity' $\tau_1(B)$, an upper bound resembling that in (7) can be obtained for $\tau_1(A^{\#})$ to yield the same qualitative

conclustion as Meyer's. Such a question is now only of interest if this can be done simply, since the right-hand side of (7) required considerable technical ingenuity at some length. We are able to obtain quite quickly (a proof is provided in the next section) the upper bound

$$\tau_1(A^\#) \leqslant \sum_{j=2}^n \frac{1}{1-\lambda_j} = \text{tr}(A^\#).$$
(10)

Since $\tau_1(A^{\#}) \ge 0$ (by (5)) it follows from (9) and (10) by the triangle inequality that

$$\frac{1}{\min_{2 \le j \le n} |1 - \lambda_j|} \le \tau_1(A^{\#}) \le \frac{n}{\min_{2 \le j \le n} |1 - \lambda_j|} \tag{11}$$

in parallel to (7). We shall see by examples that neither of the upper bounds (8) or (10) provides a particularly useful *quantitative* bound nor are they meant to. Bounds such as

$$\tau_1(A^*) \leq (1 - \tau_1(P))^{-1}$$

(Seneta, 1988, 1991) are much better quantitatively, and easily calculated from P by (5). (Indeed $0 \le \tau_1(P) \le 1$, and there are devices for coping with the situation where $\tau_1(P) = 1$.) On the other hand, the right-hand side of (10), in contrast to that of (7), shows that the situation noted by Meyer where no single λ_j is very close to 1, but enough are within range of 1 to force the right-hand side in (7) to be large does not affect the right-hand side of (10) in the same way, and confirms that such P are not badly conditioned.

2. Proof of (10)

Let $f_{ij}^{(k)}$, $k \ge 0$, be the probability of first passage from i to j in k steps ($f_{ij}^{(0)} = 0$), put $p_{ij}^{(0)} = 1$ if i = j, j = 0 if $j \ne j$, and write for $j \ne j$, where $j \ne j$ is $j \ne j$, and write for $j \ne j$, and write for $j \ne j$, and write for $j \ne j$, where $j \ne j$ is $j \ne j$, and write for $j \ne j$, where $j \ne j$ is $j \ne j$, and write for $j \ne j$ is $j \ne j$.

$$F_{ij}(z) = \sum_{k=0}^{\infty} f_{ij}^{(k)} z^k, \qquad P_{ij}(z) = \sum_{k=0}^{\infty} p_{ij}^{(k)} z^k.$$

Then it is clear that $F_{ij}(1) \le 1$, and well-known (e.g. Seneta, 1981, Section 5.4) that $P_{ii}(z) = (1 - F_{ii}(z))^{-1}$, $P_{ij}(z) = F_{ij}(z)P_{jj}(z)$, $i \ne j$. Thus for $0 < \beta < 1$, for all $i, j, P_{ij}(\beta) \le P_{jj}(\beta)$ whence

$$\sum_{k=0}^{\infty} \beta^{k} \left(p_{ij}^{(k)} - \pi_{j} \right) \leqslant \sum_{k=0}^{\infty} \beta^{k} \left(p_{jj}^{(k)} - \pi_{j} \right). \tag{12}$$

Now, according to Blackwell (1962), for our present structure of matrix P,

$$Z = I + \lim_{\beta \uparrow 1} \sum_{k=1}^{\infty} \beta^{k} (P^{k} - \mathbf{1}\boldsymbol{\pi}^{\mathrm{T}}),$$

SO

$$A^{\#} = \lim_{\beta \uparrow 1} \sum_{k=0}^{\infty} \beta^{k} (P^{k} - \mathbf{1}\boldsymbol{\pi}^{\mathrm{T}}).$$

Using (12),

$$a_{ij}^{\#} \leqslant a_{jj}^{\#}. \tag{13}$$

Since $\frac{1}{2}|a-b| = \max(a, b) - \frac{1}{2}(a+b)$,

$$\tau_1(A^{\#}) = \max_{i,j} \frac{1}{2} \sum_{s=1}^{n} \left| a_{is}^{\#} - a_{js}^{\#} \right| = \max_{i,j} \sum_{s=1}^{n} \left\{ \max(a_{is}^{\#}, a_{js}^{\#}) - \frac{1}{2} (a_{is}^{\#} + a_{js}^{\#}) \right\}$$

(repeating some lines from Seneta, 1981, p. 63). Since $A^{\#}$ has zero row sums,

$$\tau_{1}(A^{\#}) = \max_{i,j} \sum_{s=1}^{n} \max(a_{is}^{\#}, a_{js}^{\#})$$

$$\leq \sum_{s=1}^{n} a_{ss}^{\#} \text{ (by (13))}$$

$$= \operatorname{tr}(A^{\#}) \text{ (as required).}$$

An alternative direct, but longer, proof of (13) is possible.

3. Examples

Example 1. For the 8×8 matrix P in Funderlic and Meyer (1986) (the Whittaker example), the largest non-unit eigenvalue is 0.911387, $\tau_1(P) = 0.912$, $\boldsymbol{\pi}^T = (0.137, 0.049, 0.011, 0.014, 0.008, 0.050, 0.494, 0.238)$ and the matrix $A^\#$ (misprinted in Funderlic and Meyer) is

$$A^{\#} = \begin{bmatrix} 3.276 & 1.003 & -0.015 & -0.059 & 0.030 & -0.209 & -3.952 & -0.074 \\ -0.329 & 2.943 & 0.005 & -0.035 & 0.275 & -0.121 & -3.084 & 0.344 \\ -0.156 & -0.211 & 1.019 & 1.191 & 0.034 & -0.057 & -2.462 & 0.643 \\ -0.299 & -0.262 & 0.007 & 2.989 & 0.253 & -0.110 & -2.976 & 0.396 \\ -1.392 & -0.648 & -0.082 & -0.139 & 11.177 & -0.510 & -6.909 & -1.497 \\ -0.360 & -0.283 & 0.002 & -0.038 & -0.128 & 3.714 & -3.196 & 0.290 \\ -0.888 & -0.470 & -0.041 & -0.090 & -0.158 & -0.326 & 2.597 & -0.625 \\ 0.167 & -0.097 & 0.045 & 0.014 & -0.098 & 0.061 & -1.297 & 1.204 \end{bmatrix}.$$

Thus

$$\kappa(A) = 11.18 \le 2(n-1)\delta \text{ det } Z = 151.89,$$

$$\tau_1(A^\#) = 11.34 \le \text{tr}(A^\#) = 28.92,$$

$$\tau_1(A^\#) = 11.34 \le (1-\tau_1(P))^{-1} = 11.36,$$

clearly demonstrating the effectiveness of the simple bound $(1 - \tau_1(P))^{-1}$. Note also that the left-hand side of (9) for this example is 11.285, so $\kappa(A)$ does not always bound the left-side of (9), and hence (7) cannot be sharpened to the same extent.

Example 2. The intention here is to construct a well-conditioned P for which none of the non-unit eigenvalues is close to unity, but their cumulative effect makes the right-hand side of (7) large, as

envisaged by Meyer (1992),

P, aside from being bordered on the left and top, is block diagonal with $k \ 2 \times 2$ blocks. Here n = 2k + 1 is taken large, a is taken nearly unity and $\varepsilon \ll 1$. We now show that the upper bound in (7) will diverge much faster than that in (10) or that below (11). Each of the k diagonal blocks has eigenvalues $\pm a$, so P has, apart from a unit eigenvalue, k eigenvalues almost a, and k eigenvalues almost a (letting a). Also, a = 1, since a₁₁ = a₁₁ = a₁₁ = 1, a₁₂ = 1, a₁₁ = 2,..., a₁₁ (a = 1). Therefore,

$$2(n-1)\delta \prod_{j=2}^{n} \frac{1}{1-\lambda_{j}} \simeq 2(n-1)/(1-a^{2})^{(n-1)/2}$$

while

$$\sum_{j=2}^{n} \frac{1}{1-\lambda_{j}} \simeq (n-1)/(1-a^{2}).$$

On the other hand, $\tau_1(P) = a$, so

$$\tau_1(A^\#) \le (1 - \tau_1(P))^{-1} = 1/(1 - a).$$

Note that the case $\varepsilon = 0$ is covered by our results; then $\pi^T = (1, 0, 0, ..., 0)$, the irreducible sub-set consisting of the first index only. When $\varepsilon > 0$, the matrix is irreducible and aperiodic (primitive).

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