

## TWO CHARACTERIZATIONS OF INVERSE-POSITIVE MATRICES: THE HAWKINS-SIMON CONDITION AND THE LE CHATELIER-BRAUN PRINCIPLE\*

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Dedicated to the late Professors David Hawkins and Hukukane Nikaido

**Abstract.** It is shown that (a weak version of) the Hawkins-Simon condition is satisfied by any real square matrix which is inverse-positive after a suitable permutation of columns or rows. One more characterization of inverse-positive matrices is given concerning the Le Chatelier-Braun principle. The proofs are all simple and elementary.

**Key words.** Hawkins-Simon condition, Inverse-positivity, Le Chatelier-Braun principle.

**AMS subject classifications.** 15A15, 15A48.

**1. Introduction.** In economics as well as other sciences, the inverse-positivity of real square matrices has been an important topic. The Hawkins-Simon condition [9], so called in economics, requires that every principal minor be positive, and they showed the condition to be necessary and sufficient for a  $Z$ -matrix (a matrix with nonpositive off-diagonal elements) to be inverse-positive. One decade earlier, this was used by Ostrowski [12] to define an  $M$ -matrix (an inverse-positive  $Z$ -matrix), and was shown to be equivalent to some of other conditions; see Berman and Plemmons [1, Ch.6] for many equivalent conditions. Georgescu-Roegen [8] argued that for a  $Z$ -matrix it is sufficient to have only *leading* (upper left corner) principal minors positive, which was also proved in Fiedler and Ptak [5]. Nikaido's two books, [10] and [11], contain a proof based on mathematical induction. Dasgupta [3] gave another proof using an economic interpretation of indirect input.

In this paper, the Hawkins-Simon condition is defined to be the one which requires that all the *leading* principal minors should be positive, and we shall refer to it as the *weak* Hawkins-Simon condition (WHS for short). We prove that the WHS condition is necessary for a real square matrix to be inverse-positive after a suitable permutation of columns (or rows). The proof is easy and simple and uses the Gaussian elimination method. One more characterization of inverse-positive matrices is given: Each element of the inverse of the leading  $(n-1) \times (n-1)$  principal submatrix is less than or equal to the corresponding element in the inverse of the original matrix. This property is related to the Le Chatelier-Braun principle in thermodynamics.

Section 2 explains our notation, then in section 3 we present our theorems and their proofs, finally giving some numerical examples and remarks in section 4.

**2. Notation.** The symbol  $\mathbb{R}^n$  means the real Euclidean space of dimension  $n$  ( $n \geq 2$ ), and  $\mathbb{R}_+^n$  the non-negative orthant of  $\mathbb{R}^n$ . A given real  $n \times n$  matrix  $A$  is a

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map from  $\mathbb{R}^n$  into itself. The  $(i, j)$  entry of  $A$  is denoted by  $a_{ij}$ ,  $x \in \mathbb{R}^n$  stands for a column vector, and  $x_i$  denotes the  $i$ -th element of  $x$ . The symbol  $(A)_{*,j}$  means the  $j$ -th column of  $A$ , and  $(A)_{i,*}$  means the  $i$ -th row. We also use the symbol  $x_{(i)}$ , which represents the column vector in  $\mathbb{R}^{n-1}$  formed by deleting the  $i$ -th element from  $x$ . Similarly, the symbol  $A_{(i,j)}$  means the  $(n-1) \times (n-1)$  matrix obtained by deleting the  $i$ -th row and the  $j$ -th column from  $A$ . Likewise,  $A_{(:,j)}$  shows the  $n \times (n-1)$  matrix obtained by deleting the  $j$ -th column from  $A$ . The symbol  $(A)_{i,*(n)}$  shall denote the row vector formed by deleting the  $n$ -th element from  $(A)_{i,*}$ , and  $(A)_{*(n),j}$  is the column vector in  $\mathbb{R}^{n-1}$  formed by deleting the  $n$ -th element from  $(A)_{*,j}$ . The symbol  $e_i \in \mathbb{R}_+^n$  denotes a column vector whose  $i$ -th element is unity with all the remaining entries being zero.  $|A|$  denotes the determinant of  $A$ .

The inequality signs for vector comparison are as follows:

$$\begin{aligned} x &\geq y \quad \text{iff } x - y \in \mathbb{R}_+^n; \\ x &> y \quad \text{iff } x - y \in \mathbb{R}_+^n - \{0\}; \\ x &\gg y \quad \text{iff } x - y \in \text{int}(\mathbb{R}_+^n), \end{aligned}$$

where  $\text{int}(\mathbb{R}_+^n)$  means the interior of  $\mathbb{R}_+^n$ . These inequality signs are applied to matrices in a similar way.

**3. Propositions.** Let us first note that the condition “ $A$  is inverse-positive” is equivalent to the following property:

**Property 1.** For any  $b \in \text{int}(\mathbb{R}_+^n)$ , the equation  $Ax = b$  has a solution  $x \in \text{int}(\mathbb{R}_+^n)$ .

This property was used in Dasgupta and Sinha [4] to establish the nonsubstitution theorem, and in Bidard [2].

Now we can prove the following theorem.

**THEOREM 3.1.** *Let  $A$  be inverse-positive. Then the WHS condition is satisfied when a suitable permutation of columns (or rows) is made.*

*Proof.* The outline of our proof is as follows. We eliminate, step by step, a variable whose coefficient is positive. The existence of such a variable is guaranteed at each step by Property 1 above. By performing a suitable permutation of columns if necessary, this coefficient can be shown to be positively proportional to a leading principal minor of  $A$ .

Because of Property 1 above, there should be at least one positive entry in the first row of  $A$ . So, such a column and the first column can be exchanged. We assume the two columns have been permuted so that

$$a_{11} > 0.$$

Next at the second step, we divide the first equation of the system  $Ax = b$  by  $a_{11}$  and subtract this equation side by side from the  $i$ -th ( $i \geq 2$ ) equation after multiplying this by  $a_{i1}$ , to obtain

$$\begin{bmatrix} 1 & a_{12}/a_{11} & \cdots & a_{1n}/a_{11} \\ 0 & a_{22} - a_{12}a_{21}/a_{11} & \cdots & a_{2n} - a_{1n}a_{21}/a_{11} \\ \vdots & \vdots & \ddots & \vdots \\ 0 & a_{n2} - a_{12}a_{n1}/a_{11} & \cdots & a_{nn} - a_{1n}a_{n1}/a_{11} \end{bmatrix} \cdot \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{bmatrix} = \begin{bmatrix} b_1/a_{11} \\ b_2 - b_1a_{21}/a_{11} \\ \vdots \\ b_n - b_1a_{n1}/a_{11} \end{bmatrix}.$$

Notice that the  $(2, 2)$ -entry of the coefficient matrix above is

$$\frac{\begin{vmatrix} a_{11} & a_{12} \\ a_{21} & a_{22} \end{vmatrix}}{a_{11}},$$

and the corresponding entry on the RHS is

$$\frac{\begin{vmatrix} a_{11} & b_1 \\ a_{21} & b_2 \end{vmatrix}}{a_{11}}.$$

We continue this type of elimination up to the  $k$ -th step, having at the  $(k, k)$ -position

$$\frac{\begin{vmatrix} a_{11} & \cdots & \cdots & a_{1,k} \\ \vdots & \ddots & & \vdots \\ a_{k,1} & \cdots & \cdots & a_{k,k} \end{vmatrix}}{\begin{vmatrix} a_{11} & \cdots & a_{1,k-1} \\ \vdots & \ddots & \vdots \\ a_{k-1,1} & \cdots & a_{k-1,k-1} \end{vmatrix}},$$

and the RHS of the  $k$ -th equation is given as

$$\frac{\begin{vmatrix} a_{11} & \cdots & a_{1,k-1} & b_1 \\ \vdots & \ddots & & \vdots \\ a_{k,1} & \cdots & a_{k,k-1} & b_k \end{vmatrix}}{\begin{vmatrix} a_{11} & \cdots & a_{1,k-1} \\ \vdots & \ddots & \vdots \\ a_{k-1,1} & \cdots & a_{k-1,k-1} \end{vmatrix}}.$$

The denominator of these equations is known to be positive at the  $(k-1)$ -th step, and when  $b_k$  is large enough, the RHS of the  $k$ -th equation becomes positive. Thus, by Property 1, there is at least one positive coefficient in the  $k$ -th equation. Again, we assume a suitable permutation has been made so that the  $(k, k)$ -position is positive, giving

$$\begin{vmatrix} a_{11} & \cdots & \cdots & a_{1,k} \\ \vdots & \ddots & & \vdots \\ a_{k,1} & \cdots & \cdots & a_{k,k} \end{vmatrix} > 0 \quad \text{for } k = 2, 3, \dots, n.$$

Therefore, our theorem is proved for a permutation of columns. A similar result can be obtained by a suitable permutation of rows - just transpose the given matrix and apply the same proof.  $\square$

**COROLLARY 3.2.** *When  $A$  is a  $Z$ -matrix, the WHS condition is necessary and sufficient for  $A$  to be inverse-positive.*

*Proof.* First we show the necessity. Let us consider the elimination method used in the proof of Theorem 3.1. When  $A$  is a  $Z$ -matrix it is easy to notice that as elimination proceeds, a positive entry is always given at the upper left corner with the other entries (or coefficients) on the top equation being all non-positive, while the RHS of each equation always remains positive. This implies that the WHS condition holds (without any permutation).

Next we show the sufficiency. We assume that  $b \gg 0$ . When  $A$  is a  $Z$ -matrix, as elimination proceeds, a positive coefficient can appear only at the upper left corner with the remaining coefficients being all non-positive, while the RHS of each equation is always positive. So, finally we reach the equation of a single variable  $x_n$  with the two coefficients on both sides being positive. Thus,  $x_n > 0$ . Now moving backward, we find  $x \gg 0$ . Since  $b \gg 0$  is arbitrary, this proves that  $A$  is inverse-positive.  $\square$

This corollary is well known and the reader is referred to Nikaido [10, p.90, Theorem 6.1], Nikaido [11, p.14, Theorem 3.1], or Berman and Plemmons [1, p.134]. (In the diagram of Berman and Plemmons [1, p.134], the N conditions (inverse-positivity) are not connected with the E conditions (WHS) for general matrices.)

Next, we present a theorem which is related to the Le Chatelier-Braun principle; see Fujimoto [6]. This theorem is valid for a class of matrices which is more general than that of inverse-positive matrices.

**THEOREM 3.3.** *Suppose that the inverse of  $A$  has its last column and the bottom row non-negative, and that  $|A_{(n,n)}| > 0$ . Then each element of the inverse of  $A_{(n,n)}$  is less than or equal to the corresponding element of the inverse of  $A$ .*

*Proof.* It is clear that  $|A| > 0$ . The first column of the inverse of  $A$  can be obtained as a solution vector  $x \in \mathbb{R}^n$  to the system of equations  $Ax = e_1$ , while the first column of the inverse of  $A_{(n,n)}$  is a solution vector  $y \in \mathbb{R}^{n-1}$  to the system  $A_{(n,n)}y = e_{1(n)}$ . Adding these two systems with some manipulations, we get the following system:

$$(3.1) \quad A \begin{bmatrix} x_1 + y_1 \\ \vdots \\ x_{n-1} + y_{n-1} \\ x_n \end{bmatrix} = d \equiv \begin{bmatrix} 2 \\ 0 \\ 0 \\ (A)_{n,*(n)} \cdot y \end{bmatrix}.$$

By Cramer's rule, it follows that

$$x_n = \frac{|A_{(n)} d|}{|A|} = 2x_n + \frac{|A_{(n,n)}|}{|A|} \cdot (A)_{n,*(n)} \cdot y.$$

Thus, if  $x_n = (A^{-1})_{n1} > 0$ , then  $(A)_{n,*(n)} \cdot y < 0$ , and if  $x_n = 0$ , then  $(A)_{n,*(n)} \cdot y = 0$ , because  $\frac{|A_{(n,n)}|}{|A|} > 0$ .

For the  $i$ -th ( $i < n$ ) equation of (3.1), Cramer's rule gives us

$$x_i + y_i = 2x_i + \frac{|A_{(n,i)}|}{|A|} \cdot (A)_{n,*(n)} \cdot y.$$

From this, we have

$$y_i = x_i + (A^{-1})_{in} \cdot (A)_{n,*(n)} \cdot y.$$

Therefore we can assert

$$\begin{cases} y_i < x_i & \text{when } (A^{-1})_{n1} > 0 \text{ and } (A^{-1})_{in} > 0, \\ y_i = x_i & \text{when } (A^{-1})_{n1} = 0 \text{ or } (A^{-1})_{in} = 0. \end{cases}$$

For the other columns, we can proceed in a similar way by replacing  $e_1$  with the appropriate  $e_i$ .  $\square$

As a special case, we have

**COROLLARY 3.4.** *Suppose that  $A$  is inverse-positive, and the WHS condition is satisfied. Then each element of the inverse of  $A_{(n,n)}$  is less than or equal to the corresponding element of the inverse of  $A$ .*

**4. Numerical Examples and Remarks.** The first example is given by

$$A = \begin{bmatrix} -2 & 1 \\ 7 & -3 \end{bmatrix} \text{ and } A^{-1} = \begin{bmatrix} 3 & 1 \\ 7 & 2 \end{bmatrix}.$$

By exchanging two columns, we have the  $M$ -matrix

$$\begin{bmatrix} 1 & -2 \\ -3 & 7 \end{bmatrix}, \text{ whose inverse is } \begin{bmatrix} 7 & 2 \\ 3 & 1 \end{bmatrix}.$$

This satisfies the normal Hawkins-Simon condition. The inverse of (1) is (1), and the entry 1 is smaller than 7, thus verifying Corollary 3.4.

The second example is not an  $M$ -matrix:

$$A = \begin{bmatrix} 1 & -9 & 8 \\ 0 & 12 & -12 \\ -1 & 6 & -4 \end{bmatrix} \text{ and } A^{-1} = \begin{bmatrix} 2 & 1 & 1 \\ 1 & \frac{1}{3} & 1 \\ 1 & \frac{1}{4} & 1 \end{bmatrix}.$$

It should be noted that there does not exist a permutation matrix  $P$  such that  $PA$  or  $AP$  satisfies the normal Hawkins-Simon condition. However, the WHS condition is satisfied by  $A$ . The inverse of  $A_{(3,3)}$  is calculated as

$$\begin{bmatrix} 1 & -9 \\ 0 & 12 \end{bmatrix}^{-1} = \begin{bmatrix} 1 & \frac{3}{4} \\ 0 & \frac{1}{12} \end{bmatrix}.$$

This verifies Corollary 3.4.

The next example is again not an  $M$ -matrix:

$$A = \begin{bmatrix} 1 & -1 & 1 \\ 1 & 1 & -1 \\ -1 & 1 & 1 \end{bmatrix} \text{ and } A^{-1} = \begin{bmatrix} \frac{1}{2} & \frac{1}{2} & 0 \\ 0 & \frac{1}{2} & \frac{1}{2} \\ \frac{1}{2} & 0 & \frac{1}{2} \end{bmatrix}.$$

The inverse of  $A_{(3,3)}$  is calculated as

$$\begin{bmatrix} 1 & -1 \\ 1 & 1 \end{bmatrix}^{-1} = \begin{bmatrix} \frac{1}{2} & \frac{1}{2} \\ -\frac{1}{2} & \frac{1}{2} \end{bmatrix}.$$

The elements  $(A^{-1})_{11}$ ,  $(A^{-1})_{12}$ , and  $(A^{-1})_{22}$  are all equal to  $(A_{(3,3)}^{-1})_{11}$ ,  $(A_{(3,3)}^{-1})_{12}$ , and  $(A_{(3,3)}^{-1})_{22}$  because  $(A^{-1})_{32} = 0$  and  $(A^{-1})_{13} = 0$ . The entry  $(A_{(3,3)}^{-1})_{21}$  is, however,  $-\frac{1}{2}$  and is smaller than the corresponding entry  $(A^{-1})_{21} = 0$ , confirming the statements in the proof of Theorem 3.3.

The final example illustrates Theorem 3.3:

$$A = \begin{bmatrix} -\frac{17}{24} & \frac{2}{3} & -\frac{5}{24} \\ \frac{1}{6} & -\frac{1}{3} & \frac{1}{6} \\ \frac{23}{24} & -\frac{2}{3} & \frac{11}{24} \end{bmatrix} \text{ and } A^{-1} = \begin{bmatrix} -1 & -4 & 1 \\ 2 & -3 & 2 \\ 5 & 4 & 3 \end{bmatrix}.$$

Since

$$\begin{bmatrix} -\frac{17}{24} & \frac{2}{3} \\ \frac{1}{6} & -\frac{1}{3} \end{bmatrix}^{-1} = \begin{bmatrix} -\frac{8}{3} & -\frac{16}{3} \\ -\frac{4}{3} & -\frac{17}{3} \end{bmatrix},$$

these results conform to Theorem 3.3.

REMARK 4.1. The Le Chatelier-Braun principle in thermodynamics states that when an equilibrium in a closed system is perturbed, directly or indirectly, the equilibrium shifts in the direction which can attenuate the perturbation. As is explained in Fujimoto [6], the system of equations  $Ax = b$  can be solved as an optimization problem when  $A$  is an  $M$ -matrix. Thus, a solution  $x$  to the system can be viewed as a sort of equilibrium. A similar argument can be made when  $A$  is inverse-positive. That is, the solution vector  $x$  of the equations  $Ax = b$  can be obtained by solving the minimization problem:  $\min e \cdot x$  subject to  $Ax \geq b$ ,  $x \geq 0$ , where  $e$  is the row  $n$ -vector whose elements are all positive, or more simply unity. Thus, the solution vector  $x$  can be regarded as a sort of physical equilibrium. In terms of economics, the above minimization problem is to minimize the use of labor input while producing the final output vector  $b$ . (Each column of  $A$  represents a production process with a positive entry being output and a negative one input, while the vector  $e$  is the labor input coefficient vector.) Then, in our case, a perturbation is a new constraint that the  $n$ -th variable  $x_n$  should be kept constant even after the vector  $b$  shifts, destroying the  $n$ -th equation. The changes in other variables may become smaller when the increase of those variables requires  $x_n$  to be greater. This is obvious in the case of an  $M$ -matrix. What we have shown is that it is also the case with an inverse-positive matrix or even with a matrix with *positively bordered inverse* as can be seen from Theorem 3.3.

REMARK 4.2. Much more can be said about the sensitivity analysis in the case of  $M$ -matrices. We can also deal with the effects of changes in the elements of  $A$  on the solution vector  $x$ ; see Fujimoto, Herrero, and Villar [7].

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