## LEHMAN MATRICES

## GÉRARD CORNUÉJOLS, BERTRAND GUENIN, LEVENT TUNÇEL

Abstract. A pair of square 0,1 matrices $A, B$ such that $A B^{T}=E+k I$ (where $E$ is the $n \times n$ matrix of all 1 s and $k$ is a positive integer) are called Lehman matrices. These matrices figure prominently in Lehman's seminal theorem on minimally nonideal matrices. There are two choices of $k$ for which this matrix equation is known to have infinite families of solutions. When $n=k^{2}+k+1$ and $A=B$, we get point-line incidence matrices of finite projective planes, which have been widely studied in the literature. The other case occurs when $k=1$ and $n$ is arbitrary, but very little is known in this case. This paper studies this class of Lehman matrices and classifies them according to their similarity to circulant matrices.

## 1. Introduction

Let $M_{n}(K)$ denote the set of $n \times n$ matrices with elements in $K$, and let $\mathbb{B}$ denote the set $\{0,1\}$.

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We say that matrices $A, B \in M_{n}(\mathbb{B})$ form a pair of Lehman matrices if there exists a positive integer $k$ such that

$$
\begin{equation*}
A B^{T}=E+k I \tag{1}
\end{equation*}
$$

where $E$ denotes the $n \times n$ matrix of all 1 s , and $I$ is the identity matrix. Matrix $B$ is called the dual of matrix $A$. Note that $A$ is the dual of $B$ (indeed $A B^{T}=E+k I$ implies $B A^{T}=E+k I$ since $E+k I$ is symmetric). Bridges and Ryser [1] showed that every Lehman matrix is $r$-regular for some integer $r \geq 2$, i.e. it has the same number $r$ of 1 s in each row and column, see Section 2. If the dual of $A$ is $A$ itself (i.e. $A A^{T}=E+k I$ ) then $A$ is the point-line incidence matrix of a nondegenerate finite projective plane, a widely studied topic [7]. Other infinite classes of Lehman matrices occur when $k=1$ but very little is known in this case. The main purpose of this paper is to initiate a study of these matrices.

We say that $A$ is thin when $k=1$ in equation (1) and fat when $k>1$ (this terminology refers to the volume of the simplex defined by the column vectors of $A$, see Section 6.2). Nondegenerate finite projective planes with $n \geq 7$ points give rise to fat Lehman matrices. Before presenting examples of thin Lehman matrices, we introduce some notation.

Given indices $t, t^{\prime} \in[n]$ (where $[n]=\{1, \ldots, n\}$ ), a $\left(t, t^{\prime}\right)$-interval is the set of indices visited following the cyclical ordering, starting from $t$ and ending at $t^{\prime}$. We denote this interval by $\left[t, t^{\prime}\right]$. Its size is $t^{\prime}-t+1$ when $t^{\prime} \geq t$ and $t^{\prime}-t+n+1$ when $t^{\prime}<t$. Similarly, we denote the set $\{0,1, \ldots, m\}$ by $[0, m]$. Given $i \in[0, n-1]$, we say that interval $\left[t+i, t^{\prime}+i\right]$ is an $i$-shift of interval $\left[t, t^{\prime}\right]$. More generally, the $i$-shift of vector $\left(v_{1}, \ldots, v_{n}\right)$ is the vector $\left(u_{1}, \ldots, u_{n}\right)$ where $u_{j+i}=v_{j}$ if $j+i \leq n$ and $u_{j+i-n}=v_{j}$ if $j+i \geq n+1$. Vector $u$ is a shift of vector $v$ if there exists $i \in[0, n-1]$ such that $u$ is an $i$-shift of $v$.
1.1. Examples. A matrix $X \in M_{n}(\mathbb{B})$ is circulant if for all $i \in[n-1]$, row $1+i$ is an $i$-shift of row 1. Consider integers $r, s, n$ such that $r, s \geq 2$ and $r s=n+1$. We define matrices
$C_{r}^{n}, D_{s}^{n} \in M_{n}(\mathbb{B})$ as follows: $C_{r}^{n}$ and $D_{s}^{n}$ are the circulant matrices with row 1 corresponding to $[r]$ and $\{1, r, 2 r, \ldots,(s-1) r\}$ respectively. Note that $C_{r}^{n} D_{s}^{n T}=E+I$. Hence,

Remark 1.1. For all integers $r, s, n$ such that $r, s \geq 2, r s=n+1, C_{r}^{n}$ and $D_{s}^{n}$ form a pair of thin Lehman matrices.

Two matrices $X, Y$ are isomorphic if $Y$ can be obtained from $X$ by permuting the columns and the rows of $X$. If a matrix $A$ is isomorphic to a Lehman matrix, then $A$ is also a Lehman matrix (to see this, perform the same permutations on the dual and observe that (1) still holds).

2-regular Lehman matrices are perfectly understood: They are isomorphic to $C_{2}^{n}$ for $n$ odd (they are sometimes called odd holes).

Luetolf and Margot [11] enumerated all nonisomorphic Lehman matrices for $n \leq 11$. For example, they found exactly two nonisomorphic Lehman matrices for $n=8$ (to help visualize 0,1 matrices we do not write down the 0 s ):

$$
C_{3}^{8}=\left[\begin{array}{llllllll}
1 & 1 & 1 & & & & & \\
& 1 & 1 & 1 & & & & \\
& & 1 & 1 & 1 & & & \\
& & & 1 & 1 & 1 & & \\
& & & & 1 & 1 & 1 & \\
& & & & & 1 & 1 & 1 \\
1 & & & & & & 1 & 1 \\
1 & 1 & & & & & & 1
\end{array}\right] \text { and }\left[\begin{array}{llllllll}
1 & & 1 & & 1 & & & \\
& 1 & 1 & 1 & & & & \\
& & 1 & 1 & 1 & & & \\
& 1 & & 1 & & 1 & & \\
& & & & 1 & 1 & 1 & \\
& & & & & 1 & 1 & 1 \\
1 & & & & & & 1 & 1 \\
1 & 1 & & & & & & \\
& & & &
\end{array}\right] .
$$

Note that the second matrix is obtained from $C_{3}^{8}$ by adding a $0, \pm 1$ matrix of rank 1 . The main theme of this paper is that this is not a coincidence: thin Lehman matrices are either circulant matrices $C_{r}^{n}$ or "similar" to them. We make this more precise below. Define the level of a thin $r$ regular $n \times n$ Lehman matrix $A$ to be the minimum rank of $A^{\prime}-C_{r}^{n}$ over all matrices $A^{\prime}$ isomorphic to $A$. For example, the circulant matrices $C_{r}^{n}$ have level 0 and the second Lehman matrix with $n=8$ above has level 1 . To demonstrate that the notion of level is natural in the study of thin Lehman matrices, we appeal to information complexity (also known as Kolmogorov complexity).
1.2. Results. A parameter is any $\alpha \in[n]$. We say that an $n \times n$ matrix $A$ can be described with $k$ parameters $\mathcal{P}=\left\{p_{1}, \ldots, p_{k}\right\}$ if there exists an algorithm that, given $\mathcal{P}$, constructs a matrix isomorphic to $A$ (note that there is no complexity restriction on the algorithm). We prove the following theorem in Section 3.

Theorem 1.2. If $A$ is a thin $n \times n$ Lehman matrix of level $t$, then $A$ can be described with $O\left(t^{4}\right)$ parameters.

Thus thin Lehman matrices with constant level can be described with a constant number of parameters, whereas one may require $\Omega(n)$ parameters to describe a $0, \pm 1$ matrix of constant rank. This means that thin Lehman matrices with constant level are similar to $C_{r}^{n}$ in terms of information complexity.

In Section 4, we give a complete characterization of level one thin Lehman matrices, using only six parameters. This infinite class of Lehman matrices is new.

In Section 5, we prove the existence of thin Lehman matrices of arbitrarily high level and we give some constructions. In Section 6, we briefly discuss fat Lehman matrices and in Section 7 we state open problems and present some concluding remarks.
1.3. Motivation. Lehman matrices are key to understanding the set covering problem $\min \left\{c^{T} x\right.$ : $\left.M x \geq e_{m}, x \in \mathbb{B}^{n}\right\}$, a fundamental problem in combinatorial optimization (here $c$ is a given vector in $\mathbb{R}_{+}^{n}, e_{m}$ is the $m$-vector all of whose components are 1 , and $M$ is a given $m \times n$ matrix with entries equal to 0 or $1 ; x$ is the vector of unknowns). A basic question is the following: when can the set covering problem be solved by linear programming? This can be done for every objective function $c$ exactly when the set covering polytope $P:=\left\{x \in \mathbb{R}^{n}: M x \geq e_{m}, 0 \leq\right.$ $\left.x \leq e_{n}\right\}$ is integral, i.e. all its extreme points have only 0,1 components. When this occurs, the matrix $M$ is said to be ideal.

If $P$ is an integral polytope, then for all $j \in[n]$ and $\beta \in \mathbb{B}$, so are its faces $P^{\prime}:=P \cap\left\{x_{j}=\beta\right\}$. Let $P^{\prime \prime}$ be the restriction of $P^{\prime}$ to variables distinct from $x_{j}$, i.e. $P^{\prime \prime}=\left\{\left(x_{1}, \ldots, x_{j-1}, x_{j+1}, \ldots, x_{n}\right)\right.$ : $\left.\left(x_{1}, \ldots, x_{n}\right) \in P^{\prime}\right\}$. It can readily be checked that $P^{\prime \prime}$ is a set covering polytope as well, i.e. $P^{\prime \prime}=\left\{x \in \mathbb{R}^{n-1}: M^{\prime} x \geq e_{m^{\prime}}, 0 \leq x \leq e_{n-1}\right\}$ for some 0,1 matrix $M^{\prime}$. We say that $M^{\prime}$ is a minor of $M$. Thus if a matrix is ideal then so are all its minors. A 0,1 matrix is minimally nonideal if it is not ideal but all its minors are. Thus if $M$ is minimally nonideal then $P=\left\{x \in \mathbb{R}^{n}: M x \geq e_{m}, 0 \leq x \leq e_{n}\right\}$ is not an integral polytope but all the polytopes obtained from $P$ by fixing a variable $x_{j}$ to 0 or to 1 are.

An example of a minimally nonideal matrix is the point-line incidence matrix of a degenerate finite projective plane (one line contains $n-1$ points $v_{1}, \ldots, v_{n-1}$, and the remaining $n-1$ lines contain exactly two points $v_{j}, v_{n}$, for $j=1, \ldots, n-1$ ). Define the core of a minimally nonideal matrix $M$ to be the submatrix induced by those rows for which the inequalities $M \bar{x} \geq e_{m}$ hold as equality at a fractional extreme point $\bar{x}$ of $P$. Lehman [8] gave the following property of minimally nonideal matrices: If $M$ is a minimally nonideal matrix, then either it is the pointline incidence matrix of a degenerate finite projective plane or it has a unique core which is a Lehman matrix. A complete characterization of minimally nonideal matrices or of their cores seems extremely difficult. A step towards a better understanding of these matrices is to study the Lehman equation (1). This is the purpose of this paper.

A 0,1 matrix $M$ is Mengerian if for every nonnegative integral vector $c$ the linear program $\min \left\{c^{T} x: M x \geq e_{m}, 0 \leq x \leq e_{n}\right\}$ and its dual both have integral solutions. Many classical minimax theorems are associated with an underlying Mengerian matrix [3]. If a matrix is Mengerian then so are all its minors. A 0,1 matrix is minimally non-Mengerian if it is not Mengerian but all its minors are. Clearly, if $M$ is Mengerian then it is ideal. If follows that minimally non-Mengerian matrices are either minimally nonideal or ideal. In [4] it is shown that if a matrix
is minimally non-Mengerian and minimally nonideal, then its core must be thin. Hence, thin Lehman matrices are important in understanding minimally non-Mengerian matrices.

Finally, note the analogy between equation (1) and the equation $A B^{T}=E-I$ that arises in the study of perfect graphs: Lovász [10] showed that minimally imperfect graphs satisfy $A B^{T}=$ $E-I$ where $A$ ( $B$ respectively) is the maximum clique (maximum stable set respectively) versus vertex incidence matrix. Graphs that satisfy this matrix equation are called partitionable graphs and they were studied in the 1970s and following decades.

We will drop the subscript or superscript $n$ from $C_{r}^{n}, D_{s}^{n}, e_{n}$ etc. when the dimension is clear from the context.

## 2. Preliminaries

A classical result about the solutions of the Lehman matrix equation (1) was proved by Bridges and Ryser [1].

Theorem 2.1. Let $A, B \in M_{n}(\mathbb{B})$ be a Lehman pair. Then, there exist integers $r \geq 2, s \geq 2$ such that $A$ is $r$-regular, $B$ is s-regular and $r s=n+k$. Moreover, $A^{T}, B^{T}$ are also a Lehman pair.

Next, we establish that the notion of level of a Lehman matrix is invariant under duality. A matrix is 0 -regular if the sum of entries in each row and column is equal to 0 .

Proposition 2.2. Let $A, B \in M_{n}(\mathbb{B})$ be a thin Lehman pair. Then, $\operatorname{level}(A)=\operatorname{level}(B)$.

Proof. By Theorem 2.1, there exist integers $r \geq 2, s \geq 2$ such that $A$ is $r$-regular, $B$ is $s$-regular and $r s=n+1$.

Let $t=\operatorname{level}(A)$. By the definition of level, there exist $n \times n$ permutation matrices $P, Q$ such that $P A Q-C_{r}$ has rank $t$.

Claim 1. $P B Q-D_{s}$ has rankt.

Proof. We define

$$
\Sigma_{A}:=P A Q-C_{r}, \quad \Sigma_{B}:=P B Q-D_{s}
$$

Since $C_{r}$ and $D_{s}$ form a thin Lehman pair, we have

$$
\begin{aligned}
E+I & =\left(P A Q-\Sigma_{A}\right)\left(P B Q-\Sigma_{B}\right)^{T} \\
& =(P A Q)(P B Q)^{T}-C_{r} \Sigma_{B}^{T}-\Sigma_{A}(P B Q)^{T}
\end{aligned}
$$

Since $P(E+I) P^{T}=E+I$ and $A, B$ make a thin Lehman pair, so do $P A Q$ and $P B Q$. We obtain

$$
\Sigma_{B} C_{r}^{T}=-(P B Q) \Sigma_{A}^{T}
$$

By Theorem 2.1 $C_{r}^{T}$ and $D_{s}^{T}$ are a Lehman pair. Multiplying both sides of the above equation from right by $D_{s}$ and using the fact that $\Sigma_{B}$ is 0-regular, we arrive at

$$
\Sigma_{B}=-(P B Q) \Sigma_{A}^{T} D_{s}
$$

$P B Q$ and $D_{s}$ are nonsingular; therefore, $\operatorname{rank}\left(\Sigma_{B}\right)=\operatorname{rank}\left(\Sigma_{A}\right)=t$ as desired.

The above claim implies that level $(B) \leq t$. Since the roles of $A$ and $B$ are symmetric in the Lehman equation, if level $(B) \leq t-1$, we would arrive at level $(A) \leq t-1$, a contradiction. Therefore, level $(B)$ must equal $t$.

Remark 2.3. Suppose $A, B \in M_{n}(\mathbb{B})$ make a thin Lehman pair. Then using the Lehman equation,

$$
\begin{equation*}
A^{-1}=B^{T}-\frac{1}{r} E . \tag{2}
\end{equation*}
$$

Suppose $A, B$ also satisfy $A=C_{r}+\Sigma_{A}$ and $B=D_{s}+\Sigma_{B}$, where $\Sigma_{A}$ and $\Sigma_{B}$ are 0-regular matrices. Using the proof of Proposition 2.2, the identity (2), and the 0 -regularity of $\Sigma_{A}$, we
deduce

$$
\begin{equation*}
\Sigma_{B}=-B \Sigma_{A}^{T} D_{s}=-A^{-T} \Sigma_{A}^{T} D_{s}=-\left(C_{r}^{-1} \Sigma_{A} A^{-1}\right)^{T} \tag{3}
\end{equation*}
$$

## 3. INFORMATION COMPLEXITY

As we hinted in the introduction, thin Lehman matrices can be classified with respect to their relation to the circulant matrices via the notion of level. In particular, we will prove in this section that low level, thin Lehman matrices are very similar to circulant matrices. In this context, two matrices are "similar" or "close" to each other if only "little" extra information is sufficient to describe one in terms of the other. Our approach focuses on the descriptional complexity of 0,1 matrices which is in the general domain of well-known notions of Kolmogorov complexity and Shannon information theory. In such studies one has to decide ahead of time what the communicated data or the computer input "mean." (How will it be interpreted?) For our purposes, we will require that the input be treated as "positions" in an $n$-dimensional vector. While both of these areas (Kolmogorov complexity and Shannon information theory) are close to what we need, neither one is exactly suitable. Therefore, we set up our own special model below. For detailed information on Kolmogorov complexity, see [9]; for a comparison of Kolmogorov complexity and Shannon information theory, see [5].

In our approach, we are interested in describing 0,1 matrices or $0, \pm 1$ matrices. Our complexity model allows the usage of parameters in $[n]$. However, we require that any algorithm that is allowed in our model must treat these parameters as "positions" of an $n$-dimensional vector (or treat a pair of parameters as a position in an $n \times n$ matrix). For instance, to describe a 0,1 vector of length $n$, we may list the positions where contiguous ones start and end (such a representation would require $\Omega(n)$ parameters in the worst case). However, we do not allow the usage of parameters to encode the 0,1 elements as the digits of a number in $[n]$ (if this were allowed, then $\frac{n}{\log n}$ parameters would suffice to describe any 0,1 vector of length $n$ ).

As we explained in the introduction, our classification theory treats isomorphic matrices as equivalent (so does our notion of level of a thin Lehman matrix). Given thin Lehman matrices $A, A^{\prime} \in M_{n}(\mathbb{B})$, both $r$-regular, we are interested in the significant intrinsic combinatorial differences between $A$ and $A^{\prime}$. So, classification up to isomorphism also serves us well in the current section.

Let $A, B \in M_{n}(\mathbb{B})$ be a Lehman pair with $A$ being $r$-regular and $B$ being $s$-regular. To describe the 1 s in $A$, $r n$ parameters suffice. Since we allow computation (any algorithm may be used), and $A, B$ satisfy the Lehman equation, each thin Lehman matrix can be described by $\min \{r, s\} n$ parameters. E.g., if $s<r$, we describe $B$ using sn parameters and compute $A=(E+I) B^{-T}$. In contrast, one parameter suffices to describe $C_{r}$, namely $r$. Indeed, if level $(A)=O(1)$ then $O(1)$ parameters suffice to describe $A$ (see Corollary 3.8).

Given $u \in \mathbb{Z}^{n}, u_{+}, u_{-} \in \mathbb{Z}_{+}^{n}$ are the positive (negative resp.) parts of $u$ such that $u=u_{+}-u_{-}$ and $u_{+}, u_{-}$have disjoint supports. (Sometimes, we define a vector $u$ by first defining its positive and negative parts $u_{+}$and $u_{-}$and then by letting $u:=u_{+}-u_{-}$; in this latter definition, the supports of $u_{+}$and $u_{-}$need not be disjoint.) We denote the support of a vector $u$ by $\operatorname{supp}(u)$.

We say that $u \in \mathbb{Z}^{n}$ is $\left(t, C_{r}\right)$-compact if

$$
\begin{aligned}
& \operatorname{supp}\left(u_{+}\right) \subseteq \text { union of } t \text { intervals of size } r, \text { and } \\
& \operatorname{supp}\left(u_{-}\right) \subseteq \text { union of } t \text { intervals of size } r .
\end{aligned}
$$

We say that $u \in \mathbb{Z}^{n}$ is $\left(t, D_{s}\right)$-compact if

$$
\begin{aligned}
& \operatorname{supp}\left(u_{+}\right) \subseteq \text { union of the supports of } t \text { columns of } D_{s}, \text { and } \\
& \operatorname{supp}\left(u_{-}\right) \subseteq \text { union of the supports of } t \text { columns of } D_{s} .
\end{aligned}
$$

Proposition 3.1. Let $\Sigma \in M_{n}(\{0, \pm 1\})$, be 0 -regular with $\operatorname{rank}(\Sigma)=t$. If $C_{r}+\Sigma$ is nonnegative then every column and row of $\Sigma$ is $\left(t, C_{r}\right)$-compact.

Proof. We only prove that every column of $\Sigma$ is $\left(t, C_{r}\right)$-compact (our arguments directly apply to the rows of $\Sigma$ as well). First, we note that for any column $x$ of $\Sigma, x_{-}$is $\left(1, C_{r}\right)$-compact (since $C_{r}+\Sigma$ is nonnegative). Next, we prove that $x_{+}$is $\left(t, C_{r}\right)$-compact: Let $\widetilde{\Sigma}$ be the $n \times(n-1)$ matrix obtained from $\Sigma$ by deleting column $x$. Since $\Sigma$ is 0 -regular, the system:

$$
\begin{equation*}
\widetilde{\Sigma} \alpha=-x, \quad \alpha \geq 0 \tag{4}
\end{equation*}
$$

has a solution, namely $\alpha:=e$. Since $\operatorname{rank}(\widetilde{\Sigma}) \leq t$, there exists an extreme point solution $\bar{\alpha}$ of (4) such that $|\operatorname{supp}(\bar{\alpha})| \leq t$. In particular,

$$
\operatorname{supp}\left(x_{+}\right) \subseteq \bigcup_{i \in \operatorname{supp}(\bar{\alpha})} \operatorname{supp}\left(\left[\operatorname{col}_{i}(\widetilde{\Sigma})\right]_{-}\right)
$$

We conclude that $x_{+}$, and hence $x$, is $\left(t, C_{r}\right)$-compact.

Corollary 3.2. Let $\Sigma \in M_{n}(\{0, \pm 1\})$, be 0 -regular with $\operatorname{rank}(\Sigma)=t$. If $C_{r}+\Sigma$ is nonnegative then every $v \in \operatorname{rowspace}(\Sigma)$ is $\left(t^{2}, C_{r}\right)$-compact.

Proof. Choose a set of rows $\ell_{1}, \ell_{2}, \ldots, \ell_{t}$ of $\Sigma$ which forms a basis for the row space of $\Sigma$. Then $v=\sum_{i=1}^{t} \alpha_{i} \ell_{i}^{T}$, for some coefficients $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{t}$. By Proposition 3.1, each $\ell_{i}$ is $\left(t, C_{r}\right)$ compact; hence, $v$ is $\left(t^{2}, C_{r}\right)$-compact as desired.

For $p \in \mathbb{Z}^{n}$ and an $(i, j)$-interval $S \subseteq[n]$, the transition of $p$ over $S$ is

$$
\operatorname{trans}(p, S):=\sum_{k=i-1}^{j}|p(k)-p(k+1)|,
$$

where the indices are interpreted cyclically in $[n]$.
For $i, j \in[n]$, $\operatorname{dist}(i, j)$ is the size of a smallest interval containing both $i$ and $j$. Thus, if $j \geq i$, then $\operatorname{dist}(i, j)=\min \{j-i+1, i-j+n+1\}$.

Proposition 3.3. Let $r \geq 2, s \geq 2$ be integers and let $n:=r s-1$. Also let $y \in\{0, \pm 1\}^{n}$ be $\left(1, D_{s}\right)$-compact and $\ell:=C_{r}^{T} y$. Then

$$
\operatorname{trans}(\ell, S) \leq 12, \quad \text { for every interval } S \text { of size } r-1
$$

Proof. Let

$$
z_{+}:=\sum_{i \in y_{+}} \operatorname{row}_{i}\left(C_{r}\right) \text { and } z_{-}:=\sum_{i \in y_{-}} \operatorname{row}_{i}\left(C_{r}\right) .
$$

We say that $i \in[n]$ is special if $z_{+}(i) \geq 2$ or $z_{-}(i) \geq 2$. Note, $\ell=z_{+}-z_{-}$.

Claim 1. Let $i, j \in \operatorname{supp}\left(\ell_{+}\right)$be such that $\operatorname{dist}(i, j) \leq r-1$ and neither $i$ nor $j$ is special. Then $i$ and $j$ lie in the same interval of $\operatorname{supp}\left(\ell_{+}\right)$.

Proof. Clearly, $i, j \in \operatorname{supp}\left(z_{+}\right)$. Since $i, j$ are not special, $z_{+}(i)=1$ and $z_{+}(j)=1$. Let $S$ be the smallest interval containing both $i$ and $j$. Since $y$ is $\left(1, D_{s}\right)$-compact, the rows indexed by $y_{+}$ are each shifted by $r$ or $r-1$. Since $\operatorname{dist}(i, j) \leq r-1$, this implies that $S \subseteq \operatorname{supp}\left(z_{+}\right)$. Since $z_{-}(i)=z_{-}(j)=0$ and $\operatorname{dist}(i, j) \leq r-1, S \cap \operatorname{supp}\left(z_{-}\right)=\emptyset$. We conclude that $\operatorname{supp}\left(\ell_{+}\right) \supseteq S$ and that the same interval of $\ell_{+}$contains $S$.

Claim 2. There exist at most two special elements. If a special element $v$ appears in $z_{+}$then $z_{+}(v)=2$; if it appears in $z_{-}$then $z_{-}(v)=2$.

Proof. The claim follows from the matrix equation $C_{r}^{T} D_{s}=E+I$.
Let $S$ be an interval of size $r-1$. The following includes all potential contributions to $\operatorname{trans}(\ell, S)$ :

- at most 4 for each special element (by Claim 2, there are at most two such elements),
- at most 2 for each of $\ell_{+}, \ell_{-}$(by Claim 1).

The total is bounded above by 12 .

The next two remarks are useful in estimating the total number of transitions over sums of vectors and unions of intervals.

Remark 3.4. Let $\ell, \ell^{\prime} \in \mathbb{Z}^{n}$ and let $S \subseteq[n]$ be an interval. Then

$$
\operatorname{trans}\left(\ell+\ell^{\prime}, S\right) \leq \operatorname{trans}(\ell, S)+\operatorname{trans}\left(\ell^{\prime}, S\right)
$$

Remark 3.5. Let $\ell \in \mathbb{Z}^{n}$ and $S, S^{\prime} \subseteq[n]$ be intervals. Then

$$
\operatorname{trans}\left(\ell, S \cup S^{\prime}\right) \leq \operatorname{trans}(\ell, S)+\operatorname{trans}\left(\ell, S^{\prime}\right)
$$

Proposition 3.6. Let $y \in\{0, \pm 1\}^{n}$ be $\left(t, D_{s}\right)$-compact. Define $\ell:=C_{r}^{T} y$. If $\ell$ is $\left(q, C_{r}\right)$-compact, then

$$
\operatorname{trans}(\ell,[n]) \leq 48 t q
$$

Proof.

Claim 1. For every interval $S \subseteq[n]$ of size $r-1, \operatorname{trans}(\ell, S) \leq 12 t$.
Proof. Since $y$ is $\left(t, D_{s}\right)$-compact, there exist $\rho_{i} \in\{0, \pm 1\}^{n}$ such that each $\rho_{i}$ is $\left(1, D_{s}\right)$-compact and $\sum_{i=1}^{t} \rho_{i}=y$. Let $\ell_{i}:=C_{r}^{T} \rho_{i}$, for all $i \in[t]$. By Proposition 3.6, $\operatorname{trans}\left(\ell_{i}, S\right) \leq 12$. Since $\ell=\sum_{i=1}^{t} \ell_{i}$, Remark 3.4 implies the claim.

Since $\ell$ is $\left(q, C_{r}\right)$-compact,

$$
\begin{aligned}
& \operatorname{supp}\left(\ell_{+}\right) \subseteq \text { union of } q \text { intervals of size } r \\
& \operatorname{supp}\left(\ell_{-}\right) \subseteq \text { union of } q \text { intervals of size } r
\end{aligned}
$$

Therefore,

$$
\operatorname{supp}(\ell) \subseteq \text { union of } 4 q \text { intervals of size }\left\lceil\frac{r}{2}\right\rceil \leq r-1
$$

By the claim, every such interval contains at most $12 t$ transitions for $\ell$. Hence, by Remark 3.5, we have $\operatorname{trans}(\ell,[n]) \leq(4 q)(12 t)$
as desired.

Now, we are ready to prove Theorem 1.2.

Proof of Theorem 1.2. Let $A$ be an $n \times n$ thin Lehman matrix of level $t$. Then (by Theorem 2.1) $A$ is $r$-regular for some integer $r \geq 2$ and by our definition of level, there exist permutation matrices $P, Q$ such that $\operatorname{rank}\left(P A Q-C_{r}\right)=t$. Let $\Sigma_{A}:=P A Q-C_{r}$. Denote by $B$ the dual of $A$ (then $B$ is $s$-regular where $s \geq 2$ is the integer satisfying $r s=n+1$ ). Let $\Sigma_{B}:=P B Q-D_{s}$. We will describe $\Sigma_{B}$ with $O\left(t^{4}\right)$ parameters. Since the roles of $A$ and $B$ are symmetric, the same arguments also apply to $\Sigma_{A}$.

By the proof of Proposition 2.2 (or (3)), $\operatorname{rank}\left(\Sigma_{B}\right)=t$. So, there exists a $t \times t$ nonsingular submatrix $\Gamma$ of $\Sigma_{B}$ with row index set $J_{r}$, column index set $J_{c}$ such that after a suitable reordering,

$$
\Sigma_{B}=\left[\begin{array}{cc}
\Gamma & M_{1} \\
M_{2} & M_{2} \Gamma^{-1} M_{1}
\end{array}\right] .
$$

We define

$$
Y:=\left[\begin{array}{c}
\Gamma \\
M_{2}
\end{array}\right], \quad U^{T}:=\left[\begin{array}{ll}
\Gamma & M_{1}
\end{array}\right] .
$$

Further let $L:=C_{r}^{T} Y, X:=C_{r} U$. Given $L, X, J_{r}, J_{c}$ as the input, the following algorithm computes $\Sigma_{B}$ :

- Compute $D_{s}^{T} X, D_{s} L$
(this gives $\Gamma, M_{1}$ and $M_{2}$ as follows:

$$
D_{s}^{T} X=D_{s}^{T} C_{r} U=(E+I) U=U=\left[\begin{array}{c}
\Gamma^{T} \\
M_{1}^{T}
\end{array}\right]
$$

similarly,

$$
\left.D_{s} L=D_{s} C_{r}^{T} Y=(E+I) Y=Y=\left[\begin{array}{c}
\Gamma \\
M_{2}
\end{array}\right]\right)
$$

- compute $M_{2} \Gamma^{-1} M_{1}$.

We claim that $\left(L, X, J_{r}, J_{c}\right)$ can be represented by $O\left(t^{4}\right)$ parameters. Clearly, $J_{r}$ and $J_{c}$ can be represented by $t$ parameters each. So, it suffices to prove the upper bound for $L$ (since for $X$ we simply transpose the matrix $A$ ). By Corollary 3.2, every column $\ell$ of $L$ is $\left(t^{2}, C_{r}\right)$-compact. Since every column $y$ of $Y$ is a column of $\Sigma_{B}$, Proposition 3.1 implies that $y$ is $\left(t, D_{s}\right)$-compact. Now, Proposition 3.6 implies $\operatorname{trans}(\ell,[n]) \leq 48 t^{3}$. Every transition can be described by one parameter; hence, $\ell$ can be described by $O\left(t^{3}\right), L$ can be described by $O\left(t^{4}\right)$ parameters.

Remark 3.7. Theorem 1.2 also applies to partitionable matrices (those satisfying $A B^{T}=E-I$ ). We simply redefine the notion of "special" used in the proof of Proposition 3.3.

Corollary 3.8. Every pair of thin Lehman matrices with fixed level (i.e. level $(A)=t=O(1)$ ) can be described by $O(1)$ parameters.

The next section gives a complete characterization of all thin Lehman matrices of level one, using only 6 parameters.

## 4. COMPLETE CHARACTERIZATION OF LEVEL ONE MATRICES

Throughout this section $A, B \in M_{n}(\mathbb{B})$ denote level one matrices and $B$ is the dual of $A$. Moreover $A$ is $r$-regular and $B$ is $s$-regular. A matrix in $M_{n}(\mathbb{B})$ is identified with the set of pairs in $[n] \times[n]$ corresponding to its nonzero entries.

A $\left(t, q ; t^{\prime}, q^{\prime}\right)$-block is the set of pairs $(i, j)$ where $i$ is in the $\left(t, t^{\prime}\right)$-interval and $j$ is in the $\left(q, q^{\prime}\right)$-interval. A $(\rho, \sigma)$-shift of a $\left(t, q ; t^{\prime}, q^{\prime}\right)$ - block is the $\left(t+\rho, q+\sigma ; t^{\prime}+\rho, q^{\prime}+\sigma\right)$-block. A configuration $\mathcal{C}$ is a 6-tuple $\left(i, j, n_{R}, n_{C}, \rho, \sigma\right)$ associated with 4 blocks as follows. The blocks of $\mathcal{C}$ are denoted $B_{11}, B_{12}, B_{21}, B_{22}$ where $B_{11}$ is the $\left(i, j ; i+n_{R}-1, j+n_{C}-1\right)$-block, $B_{21}$ is a $(\rho, 0)$-shift of $B_{11}, B_{12}$ is a $(0, \sigma)$-shift of $B_{11}$ and $B_{22}$ is a $(\rho, \sigma)$-shift of $B_{11}$. The matrix $\Sigma(\mathcal{C})$ is defined as $-B_{11}-B_{22}+B_{21}+B_{12}$.

Theorem 4.1. A matrix $A$ is a level one (Lehman) matrix if and only if $A$ is isomorphic to $C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}$ is the configuration $\left(1,1+n_{R}, n_{R}, r-n_{R}, \operatorname{tr}, \operatorname{tr}-1\right)$ where $n_{R} \in[r-1]$ and $t \in[s-1]$.

We call any configuration of the form given in Theorem 4.1 a basic configuration. Consider, for instance, the basic configuration with $n=14, r=5, n_{R}=2, t=1$ and $\mathcal{C}=(1,3,2,3,5,4)$.

Next we describe briefly the major steps of the proof of Theorem 4.1. The "if" part is easy to check using the dual $B$ defined in Remark 4.3 below. The proof of the "only if" part consists of the following steps. Since $A$ has level one it can be written as $C_{r}+x \ell^{T}$ where $x, \ell \in\{0, \pm 1\}^{n}$. We first show in Section 4.2 that $x, \ell$ have a simple structure, i.e. only a small number of parameters are needed to describe them. This result is refined in Section 4.3 where we show that $x, \ell$ define a special type of configuration. In Section 4.4 it is proved that there exists a bijection between the configurations for $A$ and those for $B$ (after isomorphism). The proof is completed after a brief case analysis in Section 4.5.
4.1. Preliminaries. In this section, the support of a 0,1 vector $u$ will also be denoted by $u$, i.e. we use the same notation for a 0,1 vector and its support.

We say that $(P, Q)$ define the standard $\left(D_{s}, C_{s}\right)$-isomorphism if $P, Q$ are permutation matrices (of order $n$ ) such that for all indices $i, P(i,(i-1) r+1)=1$ and $Q(i, i s)=1$.

Remark 4.2. $P D_{s} Q=C_{s}$.

Proof. By definition of $D_{s}, \operatorname{row}_{i}\left(D_{s}\right)=\{i-1+t r: t \in[s]\}$. Since $Q(i, i s)=1$ and $r s=n+1$, it follows that $Q(r i, i)=1$ for all indices $i$. Now $\left(P D_{s} Q\right)_{i j}=\operatorname{row}_{i}(P) D_{s} \operatorname{col}_{j}(Q)=D_{s}((i-$ 1) $r+1, r j)$ which is equal to 1 if and only if $r j=(i-1) r+1-1+t r$ for some $t \in[s]$. We can rewrite this last condition as $r j=r(i+t)$ where $t \in[0, s-1]$. Thus $j=i+t$ where $t \in[0, s-1]$, i.e. $j \in \operatorname{row}_{i}\left(C_{s}\right)$.

We say that a permutation matrix $P$ defines a simple isomorphism if there exists $\delta \in[0, n-1]$ such that $P(i, i+\delta)=1$ for all indices $i$. Observe that $P C_{r} P^{T}=C_{r}$. Let $P, Q$ be the permutation matrices such that for all indices $i, P(i, n-i)=1$ and $Q(i, n-i+r-1)=1$. Then given $X \in M_{n}(\{0, \pm 1\}), P X Q$ is called the reverse of $X$. Note that the reverse of $C_{r}$ is $C_{r}$. Given a vector $x \in\{0, \pm 1\}^{n}$ the reverse of $x$ is $P x$. We say that $Q$ defines the standard $\left(C_{r}^{T}, C_{r}\right)$ isomorphism if $Q(i, i+r-1)=1$ for all indices $i$. Note that $C_{r}^{T} Q=C_{r}$ and that the isomorphism maps column $j$ to column $j+r-1$.

For the remainder of this section when we talk about $A$, $B$, we mean isomorphic copies $P A Q, P B Q$ such that $\operatorname{level}(A)=\operatorname{rank}\left(P A Q-C_{r}\right)$ (and by Proposition 2.2, level $(A)=$ $\left.\operatorname{level}(B)=\operatorname{rank}\left(P B Q-D_{s}\right)\right)$.

Remark 4.3. There are vectors $x, \ell, y, u \in\{0, \pm 1\}^{n}$ and $\Phi= \pm 1$ such that $A=C_{r}+x \ell^{T}$, $B=D_{s}+\Phi y u^{T}$ and $\ell=C_{r}^{T} y, x=C_{r} u, \Phi=-\frac{1}{1+x^{T} y}$. Moreover, $x^{T} e=\ell^{T} e=y^{T} e=u^{T} e=0$.

Proof. Since $A$ has level one, there exist vectors $x, \ell$ such that $A=C_{r}+x \ell^{T}$. Since $A$ is $r$-regular $x^{T} e=\ell^{T} e=0$. Define $y:=C_{r}^{-T} \ell$ and $u:=C_{r}^{-1} x$. By (2) we have $y=\left(D_{s}-\frac{1}{r} E\right) \ell=D_{s} \ell$ and $u=\left(D_{s}^{T}-\frac{1}{r} E\right) x=D_{s}^{T} x$. Moreover, $y^{T} e=\ell^{T} D_{s}^{T} e=\ell^{T} s e=0$ and similarly we can show
$u^{T} e=0$. We have,

$$
A=C_{r}+x \ell^{T}=\left(I+x \ell^{T} C_{r}^{-1}\right) C_{r}=\left(I+x y^{T}\right) C_{r} .
$$

Using Remark 6.2(1) and the above equation, we conclude that $\pm 1=\operatorname{det}\left(I+x y^{T}\right)=1+x^{T} y$. Therefore, $x^{T} y \in\{0,-2\}$ and $\Phi=-\frac{1}{1+x^{T} y}$ is well-defined and is $\pm 1$. Then it can be checked that $B=\left(I+\Phi y x^{T}\right) D_{s}$ (multiply $A B^{T}$ and use the fact that $x^{T} e=y^{T} e=0$ ). Thus

$$
B=D_{s}+\Phi y x^{T} D_{s}=D_{s}+\Phi y u^{T}
$$

Since $A$ is a 0,1 matrix, we have $x \ell^{T} \in M_{n}(\{0, \pm 1\})$. Thus, we can choose $x, \ell \in\{0, \pm 1\}^{n}$. Since $B$ is a 0,1 matrix and $\Phi= \pm 1$, we must have $y u^{T} \in M_{n}(\{0, \pm 1\})$. We established above that $y=D_{s} \ell$ and $u=D_{s}^{T} x$. Since we have $x, \ell \in\{0, \pm 1\}^{n}, y$ and $u$ are integral vectors. Therefore, $y, u \in\{0, \pm 1\}^{n}$ as desired.

Let $(P, Q)$ define the standard $\left(D_{s}, C_{s}\right)$-isomorphism. Since $B=D_{s}+\Phi y u^{T}$ it implies that $P B Q=P\left(D_{s}+\Phi y u^{T}\right) Q=P D_{s} Q+P \Phi y u^{T} Q=C_{s}+(\Phi P y)\left(Q^{T} u\right)^{T}$. Define $\tilde{y}:=\Phi P u$ and $\tilde{u}=Q^{T} u$ then $P B Q=C_{s}+\tilde{y} \tilde{u}^{T}$. Hence all results about $x, \ell$ and $A$ apply to $\tilde{y}, \tilde{u}$ and $P B Q$. The notation $\ell, x, y, u, \Phi, \tilde{y}$ and $\tilde{u}$ will be used throughout the remainder of this section.

Remark 4.4. Suppose that $\ell_{+}$is a $\left(j, j^{\prime}\right)$-interval and that $\ell_{-}$is a $\sigma$-shift of $\ell_{+}$. Suppose that $x_{-}$ is an $\left(i, i^{\prime}\right)$-interval and that $x_{+}$is a $\rho$-shift of $x_{-}$. Then we can define two distinct configurations $\mathcal{C}, \mathcal{C}^{\prime}$ from $x$ and $\ell$ such that $x \ell^{T}=\Sigma(\mathcal{C})=\Sigma\left(\mathcal{C}^{\prime}\right)$ where: $\mathcal{C}=\left(i, j, i^{\prime}-i+1, j^{\prime}-j+1, \rho, \sigma\right)$ with blocks $B_{11}=x_{-} \ell_{+}^{T}, B_{12}=x_{-} \ell_{-}^{T}, B_{21}=x_{+} \ell_{+}^{T}, B_{22}=x_{+} \ell_{-}^{T}$; and $\mathcal{C}^{\prime}=\left(i+\rho, j+\sigma, i^{\prime}-i+\right.$ $\left.1, j^{\prime}-j+1, n-\rho, n-\sigma\right)$ with blocks $B_{11}^{\prime}=x_{+} \ell_{-}^{T}, B_{12}^{\prime}=x_{+} \ell_{+}^{T}, B_{21}^{\prime}=x_{-} \ell_{-}, B_{22}^{\prime}=x_{-} \ell_{+}$.

Observe that $\mathcal{C}$ and $\mathcal{C}^{\prime}$ are determined from $x, \ell$ and the choice of $B_{11}$. Thus, we will say that $\mathcal{C}$ is the $(x, \ell)$-configuration with $B_{11}=x_{-} \ell_{+}$and that $\mathcal{C}^{\prime}$ is the $(x, \ell)$-configuration with $B_{11}=x_{+} \ell_{-}$.
4.2. $r$-structures. We use the notion of vector shift given in the introduction. A vector in $\{0, \pm 1\}^{n}$ is a type $I$, $r$-structure if it is a shift of a vector $v$ whose positive and negative parts are the intervals $v_{+}=[1, q], v_{-}=[1+t r, q+t r]$ where $q \in[r-1], t \in[s-1]$. A vector $v \in\{0, \pm 1\}^{n}$ is a Type II, $r$-structure if $v$ or $-v$ is a shift of a vector $v^{\prime}$ where $v_{+}^{\prime}=[1, q] \cup\{r\}$, $v_{-}^{\prime}=[t r, q+t r]$ and $q \in[r-2], t \in[s-1]$. A vector in $\{0, \pm 1\}^{n}$ is a Type III, $r$-structure if it is a shift of a vector $v$ where $v_{+}=[1, q] \cup\{r\}, v_{-}=[t r, q+t r-1] \cup\{q+(t-1) r\}$, where $q \in[r-2]$ and $t \in[s-1]$. The order of an $r$-structure $v$ is given by the parameter $t$. If there exists an index $\delta$ such that $\delta$ and $\delta+r-1$ are both indices of $y_{+}$(resp. $y_{-}$) then $\{\delta, \delta+r-1\}$ form a special pair of $y_{+}$(resp. $y_{-}$) and $y_{+}$(resp. $y_{-}$) is special.

Lemma 4.5. $\ell$ or its reverse is an $r$-structure of order $\left|y_{+}\right|$. Moreover, it is of type $I$ if and only if neither $y_{+}$nor $y_{-}$are special; it is of type II if and only if exactly one of $y_{+}, y_{-}$is special; it is of type III if both $y_{+}$and $y_{-}$are special.

Proof. Since $A$ does not have level $0, x_{-}, x_{+}, \ell_{-}, \ell_{+}$are all non-empty.

## Claim 1.

(1) $\ell_{+}$(resp. $\ell_{-}$) is contained in an interval of cardinality $r$.
(2) $\ell_{+}$(resp. $\ell_{-}$) is not an interval of cardinality $r$.

Proof. $A=C_{r}+x \ell^{T} \geq 0$. In particular, $C_{r}-x_{-} \ell_{+}^{T} \geq 0$, thus $\ell_{+}^{T} \subseteq \operatorname{row}_{\alpha}\left(C_{r}\right)$ for all $\alpha \in x_{-}$. This implies (1). Furthermore if $\ell_{+}$is an interval of size $r$, then $x_{-}$contains a unique element $\alpha$. Since $e^{T} x=0, x_{+}$contains a unique element $\beta$. As $C_{r}-x_{+} \ell_{-}^{T} \geq 0$, and $\ell^{T} e=0, \ell_{-}$is an interval of size $r$. Then $A$ is obtained from $C_{r}$ by permuting the rows $\alpha, \beta$, contradicting the fact that $A$ has level 1 .

## Claim 2.

(1) If $\Phi=+1$ then $y_{+} \subseteq \operatorname{col}_{\delta}\left(D_{s}\right)$, $\forall \delta \in u_{-}$. If $\Phi=-1$ then $y_{+} \subseteq \operatorname{col}_{\delta}\left(D_{s}\right)$, $\forall \delta \in u_{+}$.
(2) Consider $\delta$ such that $y_{+} \subseteq \operatorname{col}_{\delta}\left(D_{s}\right)$. Then $\forall i, j \in y_{+}, i \neq j, \operatorname{row}_{i}\left(C_{r}\right) \cap \operatorname{row}_{j}\left(C_{r}\right) \subseteq\{\delta\}$. Moreover, " $\subseteq$ " holds with " $=$ " if and only if $\{i, j\}$ is a special pair of $y_{+}$.

Proof. $D_{s}+\Phi y u^{T} \geq 0$. Suppose $\Phi=1$ as the case $\Phi=-1$ is similar. Then $D_{s}-\Phi y_{+} u_{-}^{T} \geq 0$ which implies (1). Consider $\delta$ such that $y_{+} \subseteq \operatorname{col}_{\delta}\left(D_{s}\right)$. We have $E+I=C_{r}^{T} D_{s}$ thus $e+e_{\delta}=$ $C_{r}^{T} \operatorname{col}_{\delta}\left(D_{s}\right) \geq C_{r}^{T} y_{+}=\sum_{i \in y_{+}} \operatorname{row}_{i}\left(C_{r}\right)$. Moreover, if $\delta \in \operatorname{row}_{i}\left(C_{r}\right) \cap \operatorname{row}_{j}\left(C_{r}\right)$ then $i, j$ must be a special pair. This implies (2).

We define,

$$
P:=\sum_{i \in y_{+}} \operatorname{row}_{i}\left(C_{r}\right) \quad \text { and } \quad N:=\sum_{i \in y_{-}} \operatorname{row}_{i}\left(C_{r}\right) .
$$

Then $\ell^{T}=y^{T} C_{r}=P-N$. Let $\mathcal{P}$ denote the support of $P$ and let $\mathcal{N}$ denote the support of $N$. We will show that $\mathcal{P}$ and $\mathcal{N}$ are both intervals. Partition $\mathcal{P}$ into maximal intervals $P_{1}, \ldots, P_{\alpha}$ and partition $\mathcal{N}$ into maximal intervals $N_{1}, \ldots, N_{\beta}$.

We say that sets $S, T \subseteq[n]$ cross if $S \backslash T$ and $T \backslash S$ are both non-empty.

Claim 3. $P_{i}, N_{j}$ cross for every pair $i \in[\alpha], j \in[\beta]$.

Proof. Suppose $P_{i}, N_{j}$ do not cross. We consider the case where $P_{i} \supseteq N_{j}$ as the case $P_{i} \subseteq N_{j}$ can be proved in the same way. For some indices $a, b, c, d, P_{i}=[a, b]$ and $N_{j}=[c, d]$. Since $\operatorname{row}_{a}\left(C_{r}\right)=[a, a+r-1], a \in y_{+} . \operatorname{Since} \operatorname{row}_{c}\left(C_{r}\right)=[c, c+r-1], c \in y_{-}$. As $y_{+} \cap y_{-}=\emptyset$, $a \neq c$. We omit the proof that $b \neq d$ as it is similar. Consider indices $a^{\prime}=c-1$ and $b^{\prime}=d+1$. Then $\left\{a^{\prime}, b^{\prime}\right\} \in \ell_{+}$. By Claim $1, a^{\prime}, b^{\prime}$ are contained in an interval $S$ of size at most $r$. Since $N_{j}$ is a union of rows of $C_{r},\left|N_{j}\right| \geq r$. Hence we may assume $S=\left[b^{\prime}, a^{\prime}\right]$. Since $\ell_{-} \neq \emptyset$, there exists $N_{j^{\prime}}$ where $j \neq j^{\prime} \in[\beta]$. But $N_{j^{\prime}} \subseteq S \backslash\left\{a^{\prime}, b^{\prime}\right\}$. A contradiction as $\left|N_{j^{\prime}}\right| \geq r$.

## Claim 4. $\mathcal{P}$ and $\mathcal{N}$ are both intervals.

Proof. Suppose for a contradiction $\mathcal{P}$ or $\mathcal{N}$ is not an interval. If $\mathcal{N}$ is not an interval, relabel $\ell$ by $-\ell$ and $x$ by $-x$ (as $A=C_{r}+x \ell^{T}=C_{r}+(-x)(-\ell)^{T}$ ). Then $\mathcal{P}$ becomes $\mathcal{N}$ and viceversa. Thus, we may assume there exist $P_{i_{1}}, P_{i_{2}}$ where $i_{1}, i_{2} \in[\alpha]$ and $i_{1} \neq i_{2}$. Since $\left|P_{i_{1}}\right| \geq r$, Claim 1 implies that there exists $j_{1} \in[\beta]$ such that $P_{i_{1}} \cap N_{j_{1}} \neq \emptyset$. Similarly, there exists $j_{2} \in[\beta]$ such that $P_{i_{2}} \cap N_{j_{2}} \neq \emptyset$. Note that $N_{j_{1}}, N_{j_{2}}$ need not be distinct. There exist indices $a_{1}, b_{1}, a_{2}, b_{2}, c_{1}, d_{1}, c_{2}, d_{2}$ such that $P_{i_{1}}=\left[a_{1}, b_{1}\right], P_{i_{2}}=\left[a_{2}, b_{2}\right], N_{j_{1}}=\left[c_{1}, d_{1}\right], N_{j_{2}}=\left[c_{2}, d_{2}\right]$. Since $P_{i_{1}}, N_{j_{1}}$ cross (by Claim 3) exactly one of $c_{1}, d_{1}$ is in $P_{i_{1}}$. We may assume $c_{1} \in P_{i_{1}}$ for otherwise we consider the reverse of $A$ instead of $A$, this exchanges the roles of $c_{1}$ and $d_{1}$. Since $P_{i_{2}}, N_{i_{2}}$ cross, exactly one of $c_{2}, d_{2}$ is in $P_{i_{2}}$. Thus there are two cases: (1) $c_{2} \in P_{i_{2}}$ and (2) $d_{2} \in P_{i_{2}}$.

Consider case (1). Note $N_{j_{1}} \neq N_{j_{2}}$. Then $c_{1}-1, c_{2}-1 \in \ell_{+}$. Claim 1 implies that $c_{1}-1, c_{2}-1$ are contained in an interval $S$ of cardinality $r$. But $S$ must contain strictly one of $N_{j_{1}}$ or $N_{j_{2}}$. A contradiction as $\left|N_{j_{1}}\right|,\left|N_{j_{2}}\right| \geq r$.

Consider case (2). Note $c_{1}-1, d_{2}+1 \in \ell_{+}$. Claim 1 implies $c_{1}-1, d_{2}+1$ are contained in an interval $S$ of cardinality $r$. Similarly, $b_{1}+1, a_{2}-1 \in \ell_{-}$implies $b_{1}+1, d_{2}-1$ are in an interval $S^{\prime}$ of cardinality $r$. Then $S \cup S^{\prime} \cup\left(N_{j_{1}} \backslash\left\{b_{1}+1\right\}\right) \cup\left(N_{j_{2}} \backslash\left\{a_{2}-1\right\}\right) \supseteq[n]$. Hence, $2 r+\left|N_{j_{1}}\right|+\left|N_{j_{2}}\right|-2 \geq n=r s-1$, i.e. $|\mathcal{N}| \geq\left|N_{j_{1}}\right|+\left|N_{j_{2}}\right| \geq(s-2) r+1$. It follows that $\left|y_{-}\right| \geq s-1$. If $\left|y_{+}\right|=\left|y_{-}\right|=s$ then it can be readily checked that $B$ is obtained from $D_{s}$ by permuting two columns, a contradiction as this implies $B$ (hence $A$ ) has level zero. Thus $\left|y_{+}\right|=\left|y_{-}\right|=s-1$. Claim 2(1) implies that there exists an index $\delta$ such that $y_{+} \subseteq \operatorname{col}_{\delta}\left(D_{s}\right)$. Let $i$ be the unique element in $\operatorname{col}_{\delta}\left(D_{s}\right) \backslash y_{+}$. Then $P=e+e_{\delta}-\operatorname{row}_{i}\left(C_{r}\right)$. Since $P$ decomposes into at least two intervals $P_{i_{1}}, P_{i_{2}}$, we must have $\delta \in \operatorname{row}_{i}\left(C_{r}\right)$ with $i<\delta<i+r-1$, i.e. one of the intervals $P_{i_{1}}, P_{i_{2}}$ is $\{\delta\}$. But this contradicts $\left|P_{i_{j}}\right| \geq r$ for all $j \in[\alpha]$. It follows that $\mathcal{P}$ is an interval. Similarly $\mathcal{N}$ is an interval.

By Claim 4, there are indices $a, b, c, d$ such that $\mathcal{P}=[a, b]$ and $\mathcal{N}=[c, d]$. Since $|\mathcal{P}| \geq r$, Claim 1 implies that $\mathcal{P} \cap \mathcal{N}=\emptyset$. Since by Claim 3, $\mathcal{P}$ and $\mathcal{N}$ cross, exactly one of $c, d$ is in $\mathcal{P}$. If $d \in \mathcal{P}$ then consider the reverse of $A$ instead of $A$. This will exchange the roles of $c$ and $d$, proving the result for the reverse of $\ell$. As the statement is symmetric with respect to $\ell$ and its reverse, this is acceptable. Thus, we may assume $c \in \mathcal{P}$ and $d \notin \mathcal{P}$. Let $t:=\left|y_{+}\right|=\left|y_{-}\right|$. Label elements in $y_{+}$by $\left\{i_{1}, \ldots, i_{t}\right\}$ and elements in $y_{-}$by $\left\{j_{1}, \ldots j_{t}\right\}$. We may assume that, starting from $a$ and ending at $b$, we visit rows $i_{1}, \ldots, i_{t}$ of $C_{r}$ when following the cyclic ordering. Similarly, starting from $c$ and ending at $d$, we visit rows $j_{1}, \ldots, j_{t}$ of $C_{r}$ when following the cyclic ordering.

## Claim 5.

(1) if $y_{+}$is special then $\operatorname{row}_{i_{1}}\left(C_{r}\right) \cap \operatorname{row}_{i_{2}}\left(C_{r}\right)=\{a+r-1\}$,
(2) if $y_{-}$is special then $\operatorname{row}_{j_{t-1}}\left(C_{r}\right) \cap \operatorname{row}_{j_{t}}\left(C_{r}\right)=\{d-r+1\}$.

Proof. Suppose $y_{+}$is special. Claim 2(2) implies that for some $p \in[t-1]{\text {, } \operatorname{row}_{i_{p}}\left(C_{r}\right) \cap}$ $\operatorname{row}_{i_{p+1}}\left(C_{r}\right) \neq \emptyset$. The unique element common to these rows is $a+r p-1$. Since $P(a+r p-1)=$ 2, $a+r p-1 \in \ell_{+}$. Since $a \in \ell_{+}$, Claim 1 implies that $\{a, a+r p-1\}$ is contained in an interval $S$ of size $r$. Thus $S$ does not contain $\operatorname{row}_{i_{p+1}}\left(C_{r}\right)$. It follows that $a \in \operatorname{row}_{i_{p}}\left(C_{r}\right)$, i.e. $p=1$. Then clearly $\operatorname{row}_{i_{1}}\left(C_{r}\right) \cap \operatorname{row}_{i_{2}}\left(C_{r}\right)=\{a+r-1\}$. This proves (1). The proof for (2) can be obtained by considering the reverse of $A$.

Since $r$-structures are invariant under shifting we may assume $a=1$. Let $q:=c-1$, then $[1, q] \subseteq \ell_{+}$and $[b+1, d] \subseteq \ell_{-}$.

In the remainder of the proof we consider cases depending on whether $y_{+}$and $y_{-}$are special.

Case 1. Neither $y_{+}$nor $y_{-}$are special.

Then $\ell_{+}=[1, q]$ and $\ell_{-}=[b+1, d]$. Since rows $i_{1}, \ldots, i_{t}$ of $C_{r}$ are disjoint, $b+1=1+\operatorname{tr}$ and $d=q+t r$. Hence, $\ell_{-}=[1+t r, q+t r]$ and $\ell$ is a type $\mathrm{I}, r$-structure.

Case 2. Both $y_{+}, y_{-}$are special.

By Claim 5, $\operatorname{row}_{i_{1}}\left(C_{r}\right) \cap \operatorname{row}_{i_{2}}\left(C_{r}\right)=\{r\}$, and $\operatorname{row}_{j_{t-1}}\left(C_{r}\right) \cap \operatorname{row}_{j_{t}}\left(C_{r}\right)=\{d-r+1\}$. Then $\ell_{+}=[1, q] \cup\{r\}, \ell_{-}=[b+1, d] \cup\{d-r+1\}$. Note $q \in[r-1]$. But $q \neq r-1$ because of Claim 1. Since rows $i_{1}, i_{2}$ of $C_{r}$ intersect exactly in one position and since all other pairs of rows among $i_{1}, \ldots, i_{t}$ are disjoint, $b+1=1+(\operatorname{tr}-1)=\operatorname{tr}$ and $d=q+(\operatorname{tr}-1)$. Thus $\ell_{-}=[\operatorname{tr}, q+\operatorname{tr}-1] \cup\{q+(\operatorname{tr}-1)-r+1\}$ where $q+(\operatorname{tr}-1)-r+1=q+(t-1) r$. Hence $\ell$ is a type III, $r$-structure.

Case 3. $y_{+}$is special and $y_{-}$is not special.

By Claim 5, $\operatorname{row}_{i_{1}}\left(C_{r}\right) \cap \operatorname{row}_{i_{2}}\left(C_{r}\right)=\{r\}$. Then $\ell_{+}=[1, q] \cup\{r\}$, and $\ell_{-}=[b+1, d]$. By the same argument as in Case $2, b+1=\operatorname{tr}$. Then $d=q+\operatorname{tr}$ (as we must have $\left|\ell_{+}\right|=\left|\ell_{-}\right|$). Thus $\ell_{-}=[t r, q+t r]$. Hence $\ell$ is an type II, $r$-structure.

Case 4. $y_{+}$is not special and $y_{-}$is special.

We want to show $\ell$ is a type II, $r$-structure. Since if $\ell$ is a type II, $r$-structure, so is $-\ell$, we redefine $\ell$ by $-\ell$ and $x$ by $-x$. This exchanges the roles of $\mathcal{P}$ and $\mathcal{N}$. But now $d \in \mathcal{P}$ and $c \notin \mathcal{P}$, so we consider the reverse of $A$ instead of $A$. As we exchanged $y$ for $-y$ we are in Case 3 .
4.3. Block configuration. The goal of this section is to prove:

Lemma 4.6. Let $A$ be a level one matrix. Then $A=C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}$ is a configuration $\left(i, j, n_{R}, n_{C}, t r, t^{\prime} r-\delta\right)$. where $t, t^{\prime} \in[s-1]$ and $\delta \in\{0,1\}$.

Given $S \subseteq[n] \times[n]$ we define $\operatorname{val}(S)$ to be $\left|S \cap D_{s}\right|$.

Remark 4.7. $x^{T} y=-\operatorname{val}\left(x_{-} \ell_{+}^{T}\right)-\operatorname{val}\left(x_{+} \ell_{-}^{T}\right)+\operatorname{val}\left(x_{+} \ell_{+}^{T}\right)+\operatorname{val}\left(x_{-} \ell_{-}^{T}\right) \in\{0,-2\}$.

Proof. Since $\ell=C_{r}^{T} y$,

$$
\begin{aligned}
x^{T} y & =x^{T} C_{r}^{-T} \ell=x^{T}\left(D_{s}-\frac{1}{r} E\right) \ell=x^{T} D_{s} \ell \\
& =-x_{-}^{T} D_{s} \ell_{+}-x_{+}^{T} D_{s} \ell_{-}^{T}+x_{+}^{T} D_{s} \ell_{+}+x_{-}^{T} D_{s} \ell_{-} \\
& =-\operatorname{val}\left(x_{-} \ell_{+}^{T}\right)-\operatorname{val}\left(x_{+} \ell_{-}^{T}\right)+\operatorname{val}\left(x_{+} \ell_{+}^{T}\right)+\operatorname{val}\left(x_{-} \ell_{-}^{T}\right) .
\end{aligned}
$$

Remark 4.3 states that $\Phi=-\frac{1}{1+x^{T} y}= \pm 1$. Thus $x^{T} y \in\{0,-2\}$ and the result holds.
Let $S, S^{\prime} \subseteq[n] \times[n]$. We say that $S^{\prime}$ is a horizontal translation of $S$ if $S^{\prime}$ is a $(0, t r)$-shift of $S$ where $t \in[s-1]$ and $\forall(i, j) \in S$ the numbers $j, i, i+r-1, j+t r$ do not appear in that cyclical order (note these numbers need not be all distinct). We say that $S^{\prime}$ is a vertical translation of $S$ if $S^{\prime}$ is a $(t r, 0)$-shift of $S$ where $t \in[s-1]$ and $\forall(i, j) \in S$ the numbers $i, j-r+1, j, i+t r$ do not appear in that cyclical order.

Remark 4.8. If $S^{\prime}$ is a horizontal (resp. vertical) translation of $S$ then $\operatorname{val}\left(S^{\prime}\right)=\operatorname{val}(S)$.

Proof. Let $S^{\prime}$ be a horizontal translation of $S$. Then $S$ is a $(0, t r)$-shift of $S$. Then $(i, j) \in S$ if and only if $(i, j+t r) \in S^{\prime}$. Moreover, $(i, j) \in D_{s}$ if and only if $(i, j+t r) \in D_{s}$ since $\operatorname{row}_{i}\left(D_{s}\right)=\{i, i+r-1, \ldots, i+(s-1) r-1\}$. The case for vertical translations is similar.

Remark 4.9. Let $S, S^{\prime}$ be intervals. Then $S^{\prime}$ is a $t r$-shift of $S$ for some $t \in[s-1]$ if and only if $S$ is a $\left(t^{\prime} r-1\right)$-shift of $S^{\prime}$ where $t^{\prime}=s-t \in[s-1]$.

Proof. $S^{\prime}$ is a $t r$-shift of $S$ if and only if $S$ is an $(n-t r)$-shift of $S^{\prime}$ and $n-t r=r s-1-t r=$ $t^{\prime} r-1$.

Given $S \subseteq[n] \times[n]$ and $(i, j) \in[n] \times[n]$ we abbreviate $S \backslash\{(i, j)\}$ by $S \backslash(i, j)$.

Lemma 4.10. $\ell$ is not a type II, $r$-structure.

Proof. Suppose for a contradiction, $\ell$ is a type II, $r$-structure. By considering either $x$ or $\ell$ or $-x,-\ell$ and $A$ or its reverse we may assume (after a simple isomorphism) that $\ell_{+}=[1, q] \cup\{r\}$, that $\ell_{-}=[t r, q+t r]$, and that $q \in[r-2], t \in[s-1]$. Since the smallest interval containing $\ell_{+}$ has cardinality $r,\left|x_{-}\right|=\left|x_{+}\right|=1$ and $x_{-}=\{1\}$. Applying Lemma 4.5 to $A^{T}$, it follows that $x$ or its reverse is a type I , $r$-structure. Let $\chi$ be the unique element in $x_{+}$. Remark 4.9 implies that $\chi=1+t^{\prime} r-\delta$ where $t^{\prime} \in[s-1]$ and $\delta \in\{0,1\}(\delta=1$ corresponds to the case where $x$ is a Type I, $r$-structure; $\delta=0$ corresponds to the case where the reverse of $x$ is).

Claim. $t=t^{\prime}$ and $\delta=1$.

Proof. $\operatorname{row}_{\chi}\left(C_{r}\right)=\left[1+t^{\prime} r-\delta,\left(t^{\prime}+1\right) r-\delta\right]$. Since $C_{r}-x_{+} \ell_{-}^{T} \geq 0, \ell_{-}^{T} \subseteq \operatorname{row}_{\chi}\left(C_{r}\right)$. Thus, (1) $t r \geq 1+t^{\prime} r-\delta$ and (2) $q+t r \leq\left(t^{\prime}+1\right) r-\delta$. We write (2) as $t \leq t^{\prime}+1-\frac{1}{r}(\delta+q)$. Hence $t \leq t^{\prime}$. We write (1) as $t \geq t^{\prime}+\frac{1}{r}(1-\delta)$. As $t \leq t^{\prime}$ this implies $t=t^{\prime}$ and $\delta=1$.

The claim implies that $\chi=t r$. Remark 4.9 implies that $x_{-} \ell_{-}^{T}$ is a $((s-t) r, 0)$-shift of $x_{+} \ell_{-}^{T}$. It follows that $x_{-} \ell_{-}^{T}$ is a vertical translation of $x_{+}^{T} \ell_{-}^{T}$. Hence by Remark $4.7 \operatorname{val}\left(x_{+} \ell_{-}^{T}\right)=$ $\operatorname{val}\left(x_{-} \ell_{-}^{T}\right)$. Similarly $x_{-} \ell_{+}^{T} \backslash(1,1)$ is a vertical translation of $x_{+} \ell_{+}^{T} \backslash(t r, 1)$. Hence $\operatorname{val}\left(x_{-} \ell_{+}^{T} \backslash\right.$ $(1,1))=\operatorname{val}\left(x_{+} \ell_{+}^{T} \backslash(t r, 1)\right)$. Moreover, $(1,1) \in D_{s}$ and $(\operatorname{tr}, 1) \notin D_{s}$. Thus $\operatorname{val}\left(x_{+} \ell_{+}^{T}\right)=$ $\operatorname{val}\left(x_{-} \ell_{+}^{T}\right)-1$. It follows that $-\operatorname{val}\left(x_{-} \ell_{+}^{T}\right)-\operatorname{val}\left(x_{+} \ell_{-}^{T}\right)+\operatorname{val}\left(x_{+} \ell_{+}^{T}\right)+\operatorname{val}\left(x_{-} \ell_{-}^{T}\right)=-1$, a contradiction to Remark 4.7.

A simple- $C 4$ is the matrix $\Sigma(\mathcal{C})$ where $\mathcal{C}$ is the configuration $(1,1,1,1, t r,(t+1) r-1)$. A twin- $C 4$ is the matrix $x \ell^{T}$ where $\ell_{+}=\{1\} \cup\{r\}, \ell_{-}=\{\operatorname{tr}\} \cup\{(t-1) r+1\}$ and $x_{-}=\{1\}$, $x_{+}=\{(t-1) r+1\}$ where $t \in[2, s-1]$. The order of the twin- $C 4$ is given by $t$.

Remark 4.11. Suppose $A=C_{r}+\Gamma$ where $\Gamma$ is a twin- $C 4$ of order 2 , or a simple- $C 4$. Then $A$ is isomorphic to $C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}$ is a basic configuration.

Proof. By permuting columns $r$ and $r+1$ of a twin- $C 4$ of order 2 we obtain a simple- $C 4$. By permuting rows 1 and $t r+1$ of a simple twin- $C 4$ we obtain $\Sigma(\mathcal{C})$ where $\mathcal{C}=(1,2,1, r-1, \operatorname{tr}, \operatorname{tr}-$ 1).

Lemma 4.12. Suppose $\ell$ is a type III, $r$-structure. Then after a simple isomorphism $(x, \ell)$ defines a twin-C4 of order $\left|y_{+}\right| \geq 2$.

Proof. From the hypothesis we may assume $\ell_{+}=[1, q] \cup\{r\}$, and $\ell_{-}=[t r, q+t r-1] \cup\{q+$ $(t-1) r\}$. Proceeding as in the proof of Lemma 4.10 we show that $x_{-}=\{1\}$ and $x_{+}$consists of a single element $\chi$ where $\chi=1+t^{\prime} r-\delta$ where $t^{\prime} \in[s-1]$ and $\delta \in\{0,1\}$.

Claim. $t^{\prime}=t-1, \delta=0$, and $q=1$.

Proof. Since $C_{r}-x_{+} \ell_{-}^{T} \geq 0, \ell_{-}^{T} \subseteq \operatorname{row}_{\chi}\left(C_{r}\right)=\left[1+t^{\prime} r-\delta,\left(t^{\prime}+1\right) r-\delta\right]$ and the following relation must hold: $q+(t-1) r \geq 1+t^{\prime} r-\delta$ and $q+t r-1 \leq\left(t^{\prime}+1\right) r-\delta$. We can rewrite these relations as: $t-1 \geq t^{\prime}-\frac{1}{r}(q+\delta-1)$ and $t-1 \leq t^{\prime}-\frac{1}{r}(q+\delta-1)$. It follows that $t-1=t^{\prime}-\frac{1}{r}(q+\delta-1)$. Since $t$ an integer, $q+\delta-1$ is a multiple of $r$. But $1 \leq q \leq r-2$ and $0 \leq \delta \leq 1$. It follows that $q+\delta-1=0$ hence $q=1$ and $\delta=0$. Then $t^{\prime}=t-1$.

The result follows immediately from the claim.
Consider a $\left(t, q ; t^{\prime}, q^{\prime}\right)$-block $D$. We use the following notation: ${ }^{\circ} D=(t, q), D^{\urcorner}=\left(t, q^{\prime}\right),\llcorner D=$ $\left(t^{\prime}, q\right)$ and $D_{\lrcorner}=\left(t^{\prime}, q^{\prime}\right)$. We say that ${ }^{\Gamma} D, D^{\urcorner},\left\llcorner D\right.$ and $D_{\lrcorner}$are the corners of $D$.

Lemma 4.13. If $A^{T}=C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}$ is a basic configuration, then $A$ is isomorphic to $C_{r}+\Sigma\left(\mathcal{C}^{\prime}\right)$ where $\mathcal{C}^{\prime}$ is a basic configuration.

Proof. Suppose $A^{T}=C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}=\left(1,1+n_{R}, n_{R}, r-n_{R}, t r, t r-1\right)$ where $n_{R} \in$ $[r-1]$ and $t \in[s-1]$. Let $B_{11}, B_{12}, B_{21}, B_{22}$ be blocks of $\mathcal{C}$. Then the support of $\Sigma(\mathcal{C})^{T}$ can be partitioned into blocks $B_{11}^{T}, B_{12}^{T}, B_{21}^{T}, B_{22}^{T}$. Define $B_{11}^{\prime}=B_{22}^{T}, B_{12}^{\prime}=B_{12}^{T}, B_{21}^{\prime}=B_{21}^{T}$
and $B_{22}^{\prime}=B_{11}^{T} . B_{22}$ is a $(0, t r-1)$-shift of $B_{21}$ in $\Sigma(\mathcal{C})$. Remark 4.9 implies that $B_{21}$ is a $(0,(s-t) r)$-shift of $B_{22}$ in $\Sigma(\mathcal{C})$. Thus $B_{21}^{T}=B_{21}^{\prime}$ is an $((s-t) r, 0)$-shift of $B_{22}^{T}=B_{11}^{\prime}$ in $\Sigma(\mathcal{C})^{T} . B_{22}$ is a $(t r, 0)$-shift of $B_{12}$. Remark 4.9 implies that $B_{12}$ is a $((s-t) r-1,0)$-shift of $B_{22}$ in $\Sigma(\mathcal{C})$. Thus $B_{12}^{T}=B_{12}^{\prime}$ is a $(0,(s-t) r-1)$-shift of $B_{22}^{T}=B_{11}^{\prime}$ in $\Sigma(\mathcal{C})^{T}$. Block $B_{11}^{\prime}=B_{22}^{T}$ has $r-n_{R}$ rows and $n_{R}$ columns. Let $Q$ define the standard $\left(C_{r}^{T}, C_{r}\right)$-isomorphism and let $P$ define the simple isomorphism mapping row $n_{R}+t r$ to row 1 . Then $P A Q P^{T}=$ $C_{r}+P \Sigma(\mathcal{C})^{T} Q P^{T}=C_{r}+\Sigma\left(\mathcal{C}^{\prime}\right)$ where $\mathcal{C}^{\prime}=\left(1,1+\left(r-n_{R}\right), r-n_{R}, n_{R},(s-t) r,(s-t) r-1\right)$ as ${ }^{\top} B_{11}^{\prime}=\left(1,(1+t r)+(r-1)-\left(n_{R}+t r-1\right)\right)$ where $r-1$ arises from $Q$ and $-\left(n_{R}-t r-1\right)$ arises from $P$. Observe that $\mathcal{C}^{\prime}$ is basic.

Proof of Lemma 4.6. Lemma 4.5 implies that $\ell$ or its reverse is an $r$-structure. Let $Q$ be the permutation matrix which defines the standard $\left(C_{r}^{T}, C_{r}\right)$-isomorphism. Theorem 2.1 implies that $A^{T}$ is a Lehman matrix. We have $A^{T}=C_{r}^{T}+\ell x^{T}$ thus $A^{T} Q=C_{r}^{T} Q+\ell x^{T} Q=C_{r}+\ell\left(Q^{T} x\right)^{T}$. Note that $Q^{T} x$ is an $(r-1)$-shift of $x$. Lemma 4.5 implies that $x$ or the reverse of $x$ is an $r$ structure. Lemma 4.10 implies that none of $\ell, x$, or the reverse of $\ell$ or $x$ are type II, $r$-structures.

Suppose $\ell$ or its reverse is a type I, $r$-structure. Consider the case where $x$ is a type $\mathrm{I}, r$ structure. Then $x_{-}$is a $\operatorname{tr}$-shift of $x_{+}$. Let $\mathcal{C}$ be the configuration defined by $(x, \ell)$ with $B_{11}=$ $x_{+} \ell_{-}^{T}$ (see Remark 4.4). Remark 4.9 implies that $\ell_{+}$is a $\left(t^{\prime} r-\delta\right)$-shift of $\ell_{-}$where $t^{\prime} \in[s-1]$ and $\delta \in\{0,1\}$. Then $\mathcal{C}$ is as required in the statement of Lemma 4.6. Consider the case where the reverse of $x$ is a type $\mathrm{I}, r$-structure. Then $x_{+}$is a $\operatorname{tr}$-shift of $x_{-}$. Let $\mathcal{C}$ be the configuration defined by $(x, \ell)$ with $B_{11}=x_{-} \ell_{+}^{T}$. Remark 4.9 implies that $\ell_{-}$is a $\left(t^{\prime} r-\delta\right)$-shift of $\ell_{+}$where $t^{\prime} \in[s-1]$ and $\delta \in\{0,1\}$. Then $\mathcal{C}$ is as required in Lemma 4.6.

Thus one of the following holds: (1) neither $\ell$ nor its reverse is a type $\mathrm{I}, r$-structure, (2) neither $x$ nor its reverse is a type I, $r$-structure. We will show that if (1) holds then $A=C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}$ is a basic configuration. If (2) holds, then, using the same argument (applied to $A^{T}$ instead of $A, x$ instead of $\ell$, and $\ell$ instead of $x$ ) we also obtain that $A^{T}=C_{r}+\Sigma\left(\mathcal{C}^{\prime}\right)$ where $\mathcal{C}^{\prime}$ is basic. But
then Lemma 4.13 implies that $A=C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}$ is basic. Thus Theorem 4.1 holds for $A$ and so does the weaker Lemma 4.6.

Hence it suffices to consider that (1) holds. Thus $\ell$ or its reverse is a type III, $r$-structure. We can assume we are in the former case, for if we are in the latter one, it suffices to consider $-\ell$ and $-x$ instead of $\ell$ and $x$. Lemma 4.12 implies that $(x, \ell)$ defines a twin- $C 4$ of order $\left|y_{+}\right|$. Let $(P, Q)$ define the standard $\left(D_{s}, C_{s}\right)$-isomorphism. Remarks 4.2 and 4.3 imply that $P B Q=C_{s}+\tilde{y} \tilde{u}^{T}$.

Claim. $\tilde{y}_{+}$is not an interval of cardinality $\leq s-1$.

Proof. Lemma 4.5 implies that $y_{+}$is special, i.e. there exists an index $\delta$ such that $\delta, \delta+r-1 \in y_{+}$. We have $\tilde{y}=\Phi P y$ where $P(i,(i-1) r+1)=1$ for all indices $i$ or equivalently $P(s i, i)=1$ for all indices $i$. As $\delta, \delta+r-1 \in y_{+}$, Py contains elements, $s \delta, s \delta+s r-s=s \delta-s+1$. Thus the smallest interval containing $\tilde{y}_{+}$has cardinality at least $s$.

Lemma 4.5 applied to $P B Q$ and its transpose implies that $\tilde{y}, \tilde{u}$ are $s$-structures or their reverse (note the reverse of a type III $s$-structure is equal to the inverse of a type III $s$-structure). Lemma 4.10 implies that $\tilde{y}$ is not of type II. Because of the claim, $\tilde{y}$ is not of type I either. Hence $\tilde{y}$ is of type III and Lemma 4.12 implies that $(\tilde{u}, \tilde{y})$ define a twin- $C 4$ of $(P B Q)^{T}$. In particular $\left|\tilde{y}_{+}\right|=\left|y_{+}\right|=2$. Hence $(x, \ell)$ is a twin- $C 4$ of order 2. Then Remark 4.11 completes the proof.
4.4. Block configurations in the dual. The goal of this section is to prove the following result.

Lemma 4.14. Suppose $A=C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}$ is a configuration $\left(i, j, n_{R}, n_{C}, t r, t^{\prime} r-\delta\right)$ where $t, t^{\prime} \in[s-1]$ and $\delta \in\{0,1\}$. Let $(P, Q)$ define the standard $\left(D_{s}, C_{s}\right)$-isomorphism. Then $P B Q=C_{s}+\Sigma\left(\mathcal{C}^{\prime}\right)$ where $\mathcal{C}^{\prime}$ has the following parameters:
(1) If $\Phi=+1$ and $\delta=0$ then $\mathcal{C}^{\prime}=\left(\tilde{\jmath}, \tilde{\imath}, t^{\prime}, t, n_{C} s, n_{R} s\right)$,
(2) If $\Phi=-1$ and $\delta=0$ then $\mathcal{C}^{\prime}=\left(\tilde{\jmath}, \tilde{\imath}+n_{R} s, t^{\prime}, t, n_{C} s,\left(r-n_{R}\right) s-1\right)$,
(3) If $\Phi=+1$ and $\delta=1$ then $\mathcal{C}^{\prime}=\left(\tilde{\jmath}-\left(s-t^{\prime}\right), \tilde{\imath}+n_{R} s, s-t^{\prime}, t, n_{C} s,\left(r-n_{R}\right) s-1\right)$,
(4) If $\Phi=-1$ and $\delta=1$ then $\mathcal{C}^{\prime}=\left(\tilde{\jmath}-\left(s-t^{\prime}\right), \tilde{\imath}, s-t^{\prime}, t, n_{C} s, n_{R} s\right)$,
where $\tilde{\imath}=(i-1) s+1$ and $\tilde{\jmath}=(j-1) s+1$.

We will need a number of preliminary results.

Lemma 4.15. Suppose $(P, Q)$ defines the standard $\left(D_{s}, C_{s}\right)$-isomorphism. Let $v$ be an $(a, b)$ interval and $|v| \leq r-1$. Let $\tilde{a}=(a-1) s+1$ and let $\tilde{b}=b s$. Then
(1) $P D_{s} v$ is an $(\tilde{a}, \tilde{b})$-interval,
(2) $Q^{T} D_{s}^{T} v$ is an $(\tilde{a}, \tilde{b})$-interval.

Proof. Consider part (1). Note $P D_{s} Q=C_{s}$, thus $P D_{s}=C_{s} Q^{T}$ which implies that $P D_{s} v=$ $C_{s} Q^{T} v$. Since $Q(i, s i)=1, \operatorname{col}_{i}\left(C_{s} Q^{T}\right)=\operatorname{col}_{s i}\left(C_{s}\right)$. Thus $C_{s} Q^{T} v=\sum_{i \in v} \operatorname{col}_{s i}\left(C_{s}\right)$. Note $\operatorname{col}_{s i}\left(C_{s}\right)=[(i-1) s+1, i s]$. Thus for any index $i, \operatorname{col}_{s i}\left(C_{s}\right) \cap \operatorname{col}_{s(i+1)}\left(C_{s}\right)=\emptyset$ and $\operatorname{col}_{s i}\left(C_{s}\right) \cup$ $\operatorname{col}_{s(i+1)}\left(C_{s}\right)$ forms an interval. It follows that $C_{s} Q^{T} v$ is the required interval.

Consider part (2). Note $Q^{T} D_{s}^{T} P^{T}=C_{s}^{T}$, thus $Q^{T} D_{s}^{T}=C_{s}^{T} P$ which implies that $Q^{T} D_{s}^{T} v=$ $C_{s}^{T} P v$. Since $P(i,(i-1) r+1)=1$ we have that $P(i s, i-1+s)=1$ and $P((i-1) s+1, i)=1$ for all indices $i$. Hence $\operatorname{row}_{i}\left(C_{s}^{T} P\right)=\operatorname{row}_{(i-1) s+1}\left(C_{s}\right)=[(i-1)+1, i s]$ and $Q^{T} D_{s}^{T} v=C_{s}^{T} P v=$ $\sum_{i \in v} \operatorname{row}_{(i-1) s+1}\left(C_{s}\right)$. Proceed now as in part (1).

Lemma 4.16. Let $t, \Delta \in[s-1]$ and $a \in[n]$ and let $\delta \in\{0,1\}$. Suppose $\ell_{+}$is an $(a, a+\Delta-1)$ interval and $\ell_{-}$is a $(t r-\delta)$-shift of $\ell_{+}$. Define y by $\ell=C_{r}^{T}$ y and let $\tilde{a}=(a-1) s+1$.
(1) If $\delta=0$ then $(P y)_{+}=[\tilde{a}, \tilde{a}+t-1]$ and $(P y)_{-}$is a $\Delta s$-shift of $(P y)_{+}$.
(2) If $\delta=1$ then $(P y)_{-}=[\tilde{a}-(s+t), \tilde{a}-1]$ and $(P y)_{+}$is a $\Delta s$-shift of $(P y)_{-}$.
(3) Statement (1) remains true if we replace $\ell$ by $x^{\prime}, y$ by $u^{\prime}$ and Py by $Q^{T} u^{\prime}$ where $x^{\prime}=C_{r} u^{\prime}$.

Proof. Consider part (1). We have $\ell=C_{r}^{T} y$. Thus $y=C_{r}^{-T} \ell=\left(D_{s}-\frac{1}{r} E\right) \ell=D_{s} \ell$ where the last equality follows from the fact that $\ell$ is 0-regular. Define $\mathcal{P}=P D_{s} \ell_{+}$and $\mathcal{N}=P D_{s} \ell_{-}$.

Since $\ell=\ell_{+}-\ell_{-}$it follows that $P y=P D_{s} \ell_{+}-P D_{s} \ell_{-}=\mathcal{P}-\mathcal{N}$. Applying Lemma 4.15(1) we obtain that $\mathcal{P}=[\tilde{a}, \tilde{b}]$ where $\tilde{a}=(a-1) s+1$ and $\tilde{b}=(a+\Delta-1) s=\tilde{a}+\Delta s-1$. Applying Lemma 4.15(1) we also obtain that $\mathcal{N}=\left[\tilde{a}^{\prime}, \tilde{b}^{\prime}\right]$ where $\tilde{a}^{\prime}=(a+t r-1) s+1=\tilde{a}+\operatorname{tr} s=\tilde{a}+t$ and $\tilde{b}^{\prime}=(a+\Delta-1+t r) s=\tilde{b}+t=(\tilde{a}+\Delta s-1)+t$. Hence, $\mathcal{N}$ is a $t$-shift of $\mathcal{P}$. Since $t<s, \emptyset \neq \mathcal{P} \cap \mathcal{N}=[\tilde{a}+t, \tilde{b}]$. It follows that $(P y)_{+}=\mathcal{P}-\mathcal{N}=[\tilde{a}, \tilde{a}+t-1]$ and $(P y)_{-}=\mathcal{N}-\mathcal{P}=[\tilde{b}+1, \tilde{b}+t]=[\tilde{a}+\Delta s, \tilde{a}+\Delta s+t-1]$. Hence (1) holds.

Consider case (2). We define $\mathcal{P}$ and $\mathcal{N}$ in the same manner as in case (1). Applying Lemma 4.15(1) to $\mathcal{P}$ we obtain that (as in case (1)) $\mathcal{P}=[\tilde{a}, \tilde{b}]$ where $\tilde{b}=\tilde{a}+\Delta s-1$. Applying Lemma 4.15(1) to $\mathcal{N}$ we obtain that $\mathcal{N}=\left[\tilde{a}^{\prime}, \tilde{b}^{\prime}\right]$ where $\tilde{a}^{\prime}=(a+\operatorname{tr}-1-1) s+1=\tilde{a}+t-s$ and $\tilde{b}^{\prime}=(a+\Delta-1+\operatorname{tr}-1) s=\tilde{a}+\Delta s-1-(s-t)=\tilde{b}-(s-t)$. Thus $\mathcal{N}$ is a $(t-s)$-shift of $\mathcal{P}$. As $t \leq s, \mathcal{P} \cap \mathcal{N}=[\tilde{a}, \tilde{b}-(s-t)]$. It follows that $(P y)_{-}=\mathcal{N} \backslash \mathcal{P}=[\tilde{a}-(s-t), \tilde{a}-1]$ and $(P y)_{+}=\mathcal{P} \backslash \mathcal{N}=[\tilde{a}-(s-t), \tilde{a}-1]$ and $(P y)_{+}=\mathcal{P} \backslash \mathcal{N}=[\tilde{b}-(s-t)+1, \tilde{b}]=$ $[\tilde{a}+\Delta s-(s-t), \tilde{a}+\Delta s-1]$. This proves (2).

Consider case (3). We have $x^{\prime}=C_{r} u^{\prime}$, thus $u^{\prime}=C_{r}^{-1} x^{\prime}=\left(D_{s}^{T}-\frac{1}{r} E\right) x^{\prime}=D_{s}^{T} x^{\prime}$. Define $\mathcal{P}=Q^{T} D_{s}^{T} x_{+}^{\prime}$ and $\mathcal{N}=Q^{T} D_{s}^{T} x_{-}^{\prime}$. Since $x^{\prime}=x_{+}^{\prime}-x_{-}^{\prime}$ it follows that $Q^{T} u^{\prime}=Q^{T} D_{s}^{T} x_{+}^{\prime}-$ $Q^{T} D_{s}^{T} x_{-}^{\prime}=\mathcal{P}-\mathcal{N}$. Using Lemma 4.15(2) we obtain that $\mathcal{P}, \mathcal{N}$ are the same intervals that as in part (1). The proof now proceeds in the same way.

We are now ready for the main result of this section.

Proof of Lemma 4.14. We have $\mathcal{C}=\left(i, j, n_{R}, n_{C}, \operatorname{tr}, t^{\prime} r-\delta\right)$ and $\Sigma(\mathcal{C})=x \ell^{T}$ for some $x, \ell \in$ $\{0, \pm 1\}^{n}$. We can choose $x, \ell$ such that $x_{-}=\left[i, i+n_{R}-1\right], x_{+}$is a tr-shift of $x_{-} ; \ell_{+}=$ $\left[j, j+n_{C}-1\right], \ell_{-}$is a $\left(t^{\prime} r-\delta\right)$-shift of $\ell_{+}$. Recall that $\Sigma\left(\mathcal{C}^{\prime}\right)=\tilde{y} \tilde{u}^{T}$ where $\tilde{y}=\Phi P y$ and $\tilde{u}=Q^{T} u$. Let $x^{\prime}=-x$ and $u^{\prime}=-u$. Since $x=C_{r} u, x^{\prime}=C_{r} u^{\prime}$. Lemma 4.16(3) implies that $\left(Q^{T} u^{\prime}\right)_{+}=\tilde{u}_{-}=[\tilde{\imath}, \tilde{\imath}+t-1]$ and $\left(Q^{T} u^{\prime}\right)_{-}=\tilde{u}_{+}$is a $n_{R} s$-shift of $u_{-}$where $\tilde{\imath}=(i-1) s+1$.

Consider part (1), i.e $\Phi=1, \delta=0$. Then the relation $\tilde{y}=P y$ and Lemma 4.16(1) imply that $\tilde{y}_{+}=(P y)_{+}=\left[\tilde{\jmath}, \tilde{\jmath}+t^{\prime}-1\right]$ and $\tilde{y}_{-}=(P y)_{-}$is a $n_{C} s$-shift of $\tilde{y}_{+}$, where $\tilde{\jmath}=(j-1) s+1$. Let $\mathcal{C}^{\prime}$ be the configuration defined by $(\tilde{y}, \tilde{u})$ with $B_{11}^{\prime}=\tilde{y}_{+} \tilde{u}_{-}^{T}$ (see Remark 4.4). The first two parameters of $\mathcal{C}^{\prime}$ are given by the corner ${ }^{「} B_{11}^{\prime}=(\tilde{\jmath}, \tilde{\imath})$ and each of the blocks have $t^{\prime}$ rows and $t$ columns.

Consider part (2), i.e. $\Phi=-1$ and $\delta=0$. Then $\tilde{y}=-P y$ and Lemma 4.16(1) implies that $\tilde{y}_{-}=(P y)_{-}=\left[\tilde{\jmath}, \tilde{\jmath}+t^{\prime}-1\right]$ and $\tilde{y}_{+}=(P y)_{-}$is an $n_{C} s$-shift of $\tilde{y}_{-}$(and $\tilde{\jmath}$ is as above). Let $\mathcal{C}^{\prime}$ be the configuration defined by $(\tilde{y}, \tilde{u})$ with $B_{11}^{\prime}=\tilde{y}_{-} \tilde{u}_{+}^{T}$. The first two parameters of $\mathcal{C}^{\prime}$ are given by the corner ${ }^{'} B_{11}^{\prime}=\left(\tilde{\jmath}, \tilde{\imath}+n_{R} s\right)$ and each of the blocks have $t^{\prime}$ rows and $t$ columns. Since $\tilde{u}_{+}$is an $n_{R} s$-shift of $\tilde{u}_{-}$, Remark 4.9 implies that $\tilde{u}_{-}$is an $\left(\left(r-n_{R}\right) s-1\right)$-shift of $\tilde{u}_{+}$.

Consider part (3), i.e. $\Phi=1$ and $\delta=1$. Then $\tilde{y}=P y$ and Lemma 4.16(2) implies that $\tilde{y}_{-}=(P y)_{-}=\left[\tilde{\jmath}-\left(s-t^{\prime}\right), \tilde{\jmath}-1\right]$ and $\tilde{y}_{+}=(P y)_{+}$is an $n_{C} s$-shift of $\tilde{y}_{-}$. Let $\mathcal{C}^{\prime}$ be the configuration defined by $(\tilde{y}, \tilde{u})$ with $B_{11}^{\prime}=\tilde{y}_{-} \tilde{u}_{+}^{T}$. Note that ${ }^{「} B_{11}^{\prime}=\left(\tilde{\jmath}-\left(s-t^{\prime}\right), \tilde{\imath}+n_{R} s\right)$ and that the blocks have $s-t^{\prime}$ rows and $t$ columns. Since $\tilde{u}_{+}$is an $n_{R} s$-shift of $\tilde{u}_{-}$, Remark 4.9 implies that $\tilde{u}_{-}$is an $\left(\left(r-n_{R}\right) s-1\right)$-shift of $\tilde{u}_{+}$.

Consider part (4), i.e. $\Phi=-1$ and $\delta=1$. Then $\tilde{y}=-P y$ and Lemma 4.16(2) implies that $\tilde{y}_{+}=(P y)_{-}=\left[\tilde{\jmath}-\left(s-t^{\prime}\right), \tilde{\jmath}-1\right]$ and $\tilde{y}_{-}=(P y)_{+}$is an $n_{C} s$-shift of $\tilde{y}_{+}$. Let $\mathcal{C}^{\prime}$ be the configuration defined by $(\tilde{y}, \tilde{u})$ with $B_{11}^{\prime}=\tilde{y}_{+} \tilde{u}_{-}^{T}$. Note that ${ }^{\top} B_{11}=\left(\tilde{\jmath}-\left(s-t^{\prime}\right), \tilde{\imath}\right)$ and the blocks have $s-t^{\prime}$ rows and $t$ columns.
4.5. Case analysis. Lemma 4.6 implies (after possibly a simple isomorphism) that $A=C_{r}+$ $\Sigma(\mathcal{C})$ where $\mathcal{C}$ is a configuration $\left(1, b, n_{R}, n_{C}, t r, t^{\prime} r-\delta\right)$ where $b$ is an index, $n_{R}, n_{C} \in[r-$ $1], t, t^{\prime} \in[s-1]$ and $\delta \in\{0,1\}$. Let $B_{11}, B_{12}, B_{21}, B_{22}$ denote the blocks of $\mathcal{C}$. The variables $b, n_{R}, n_{C}, t, t^{\prime}$ and $\delta$ are used throughout the remainder of this section.

Lemma 4.17. We may assume $t=t^{\prime}$.

Proof. Note ${ }^{「} B_{22}=\left(1+t r, b+t^{\prime} r-\delta\right)$. Thus $b+t^{\prime} r-\delta \in \operatorname{row}_{1+t r}\left(C_{r}\right)$ i.e. there exists $q \in[0, r-1]$ such that $b+t^{\prime} r-\delta=1+t r+q$, i.e $r\left(t^{\prime}-t\right)=q-b+\delta+1$. As $b \leq r, q-b+\delta+1>-r$; hence $t^{\prime}-t \geq 0$. Suppose $t^{\prime}-t \geq 1$. Then $q-b+\delta+1$ is a multiple of $r$, but as $q \leq r-1, b \geq 1$ and $\delta \leq 1$ we must have $q=r-1, b=\delta=1$. As ${ }^{\ulcorner } B_{11}=(1, b)=(1,1), n_{R}=1$ and as ${ }^{\ulcorner } B_{22}=(1+\operatorname{tr},(t+1) r), n_{C}=1$. Thus $\mathcal{C}=(1,1,1,1, t r,(t+1) r-1)$, i.e. it is a simple- $C 4$. We are then done by Remark 4.11.

Thus throughout the remainder of the section $t=t^{\prime}$.

Lemma 4.18. If $\delta=0$ then $\operatorname{val}\left(B_{11}\right)=\operatorname{val}\left(B_{22}\right)$ and $\operatorname{val}\left(B_{11}\right) \neq 1$.

Proof. Since $B_{11} \subseteq C_{r}, B_{11} \cap D_{s} \subseteq\left\{B_{11}^{\urcorner},\left\llcorner B_{11}\right\}\right.$. Similarly, $B_{22} \cap D_{s} \subseteq\left\{B_{22}^{\urcorner},\left\llcorner B_{22}\right\}\right.$. Since $t=t^{\prime}$ and $\delta=\emptyset, B_{11}^{\urcorner} \cap D_{s} \neq \emptyset$ if and only if $B_{22}^{\urcorner} \cap D_{s} \neq \emptyset$ and $B_{11} \cap D_{s} \neq \emptyset$ if and only if $B_{22} \neq \emptyset$. It follows that $\operatorname{val}\left(B_{11}\right)=\operatorname{val}\left(B_{22}\right)$. Suppose $\operatorname{val}\left(B_{11}\right)=\operatorname{val}\left(B_{22}\right)=1$. Assume $B_{11}^{\urcorner} \in D_{s}$ as the case $B_{11} \in D_{s}$ can be dealt with similarly. Then $B_{22}^{\urcorner} \in D_{s}$. Since $\delta=0, B_{12}$ is a horizontal translation of $B_{11}$, hence Remark 4.8 implies that $\operatorname{val}\left(B_{11}\right)=\operatorname{val}\left(B_{12}\right)$. $B_{22} \backslash B_{22}^{\urcorner}$is a horizontal translation of $B_{21} \backslash B_{21}^{\urcorner}$, hence Remark 4.8 implies that val $\left(B_{21} \backslash B_{21}^{\urcorner}\right)=$ $\operatorname{val}\left(B_{22} \backslash B_{22}^{\urcorner}\right)$. Moreover, $B_{21}^{\urcorner} \notin D_{s}$ and $B_{22}^{\urcorner} \in D_{s}$. It follows that $\operatorname{val}\left(B_{21}\right)-\operatorname{val}\left(B_{22}\right)=-1$. Hence $-\operatorname{val}\left(B_{11}\right)+\operatorname{val}\left(B_{12}\right)+\operatorname{val}\left(B_{21}\right)-\operatorname{val}\left(B_{22}\right)=-1$, contradicting Remark 4.7.

Lemma 4.19. Let $(P, Q)$ define the standard $\left(C_{s}, D_{s}\right)$-isomorphism. (1) Suppose $A=C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}$ is basic, then $P B Q=C_{s}+\Sigma\left(\mathcal{C}^{\prime}\right)$ where $\mathcal{C}^{\prime}$ is basic. (2) Suppose $P B Q=C_{s}+\Sigma\left(\mathcal{C}^{\prime}\right)$ where $\mathcal{C}^{\prime}$ is basic then $A=C_{r}+\Sigma(\mathcal{C})$ where $\mathcal{C}$ is basic.

Proof. Since we can interchange the roles of $A$ and $P B Q$ it suffices to prove (1). $B_{21} \backslash B_{21}$ is a vertical translation of $B_{11} \backslash B_{11}^{\urcorner}$. Remark 4.8 implies that $\operatorname{val}\left(B_{11} \backslash B_{11}^{\urcorner}\right)=\operatorname{val}\left(B_{21} \backslash B_{21}^{\urcorner}\right)$. Moreover, $B_{11}^{\urcorner} \in D_{s}$; but $B_{21}^{\urcorner} \notin D_{s}$. Thus, $\operatorname{val}\left(B_{11}\right)=\operatorname{val}\left(B_{21}\right)+1$. Similarly, we prove that $\operatorname{val}\left(B_{22}\right)=\operatorname{val}\left(B_{12}\right)+1$. Hence $-\operatorname{val}\left(B_{11}\right)+\operatorname{val}\left(B_{12}\right)+\operatorname{val}\left(B_{21}\right)-\operatorname{val}\left(B_{22}\right)=-2$.

Remark 4.7 implies that $x^{T} y=-2$ hence (Remark 4.3) $\Phi=+1$. Thus, we are in case (3) of Lemma 4.14 with $i=1$ and $j=1+n_{R}$. Then $\tilde{\imath}=1$ and $\tilde{\jmath}=n_{R} s+1$. Thus $\mathcal{C}^{\prime}=$ $\left(n_{R} s+1-(s-t), n_{R} s+1, s-t, t,\left(r-n_{R}\right) s,\left(r-n_{R}\right) s-1\right)$. After a simple isomorphism, mapping row $n_{R} s+1-(s-t)$ to 1 , we have $\mathcal{C}^{\prime}=\left(1,(s-t)+1, s-t, t,\left(r-n_{R}\right) s,\left(r-n_{R}\right) s-1\right)$. Define $n_{R}^{\prime}=s-t$ and $q=r-n_{R}$, then $\mathcal{C}^{\prime}=\left(1,1+n_{R}^{\prime}, n_{R}^{\prime}, s-n_{R}^{\prime}, q s, q s-1\right)$ which is basic.

We can now prove the main theorem of this section.

Proof of Theorem 4.1: The "if" part of the statement follows from Lemma 4.19. Let $\mathcal{C}^{\prime}$ be the configuration obtained from $\mathcal{C}=\left(i, j, n_{R}, n_{C}, t r, t^{\prime} r-\delta\right)$ in Lemma 4.14 where $i=1$ (we will consider each of the 4 cases of the lemma separately). Note, $\tilde{\imath}=1$. Denote by $B_{11}, B_{12}, B_{21}, B_{22}$ the blocks corresponding to $\mathcal{C}$.

Case 1. $\Phi=1$ and $\delta=0$.

Then $\mathcal{C}^{\prime}=\left(\tilde{\jmath}, 1, t, t, n_{C} s, n_{R} s\right)$. By applying Lemma 4.17 to $P B Q$ instead of $A$ we obtain that $n_{C}=n_{R}$. Suppose $\operatorname{val}\left(B_{11}\right)=\operatorname{val}\left(B_{22}\right)=0$. Then $B_{12}$ is a horizontal translation of $B_{11}$ and $B_{22}$ is a horizontal translation of $B_{21}$. Remark 4.8 implies that $\operatorname{val}\left(B_{11}\right)=\operatorname{val}\left(B_{12}\right)$ and $\operatorname{val}\left(B_{21}\right)=\operatorname{val}\left(B_{22}\right)$. Then $-\operatorname{val}\left(B_{11}\right)+\operatorname{val}\left(B_{12}\right)+\operatorname{val}\left(B_{21}\right)-\operatorname{val}\left(B_{22}\right)=0$. Remark 4.7 implies that $x^{T} y=0$. Remark 4.3 implies that $\Phi=-1$, a contradiction. Lemma 4.18 implies $\operatorname{val}\left(B_{11}\right)=$ $\operatorname{val}\left(B_{22}\right) \neq 1$. Hence $\operatorname{val}\left(B_{11}\right)=\operatorname{val}\left(B_{22}\right) \geq 2$. Thus $\operatorname{val}\left(B_{11}\right)=2$ and $\left\{{ }_{L} B_{11}, B_{11}^{\urcorner}\right\} \subseteq D_{s}$. It follows that $n_{C}=n_{R}=\frac{r+1}{2}$ and that $r$ is odd. Since $i=1$ we must have $j=r-n_{C}+1=\frac{r+1}{2}$. It follows that $\tilde{\jmath}=\left(\frac{r-1}{2}\right) s+1=\frac{1}{2}(n+1-s)+1$. We must have $(\tilde{\imath}, \tilde{\jmath}) \in C_{s}$ thus $\tilde{\jmath} \in \operatorname{col}_{1}\left(C_{s}\right)=$ $\{n-s+2, \ldots, n\} \cup\{1\}$, i.e. $\frac{1}{2}(n+1-s)+1 \geq n-s+2$, which implies $1 \geq n-s+2$, a contradiction.

Case 2. $\Phi=-1$ and $\delta=0$.

Then $\mathcal{C}^{\prime}=\left(\tilde{\jmath}, 1+n_{R} s, t, t, n_{C} s,\left(r-n_{R}\right) s-1\right)$. By applying Lemma 4.17 to $P B Q$ instead of $A$ we obtain that $n_{C}=r-n_{R}$. It follows that exactly one of $B_{11}, B_{11}^{\urcorner}$is in $D_{s}$, i.e. that $\operatorname{val}\left(B_{11}\right)=1$. But this contradicts Lemma 4.18.

Case 3. $\Phi=1$ and $\delta=1$.

Then $\mathcal{C}^{\prime}=\left(\tilde{\jmath}-(s-t), 1+n_{R} s, s-t, t, n_{C} s,\left(r-n_{R}\right) s-1\right)$. By applying Lemma 4.17 to $P B Q$ instead of $A$ we obtain that $n_{C}=r-n_{R}$. Then exactly one of $B_{11}, B_{11}^{\urcorner}$is in $D_{s}$. By Lemma 4.18 exactly one of $B_{22}, B_{22}^{\urcorner}$is in $D_{s}$. Moreover, since $\delta=1$, we must have $B_{11}^{\urcorner} \in D_{s}$ and $B_{22} \in D_{s}$. Since $i=1, j=r-n_{C}+1=r-\left(r-n_{R}\right)+1=n_{R}+1$. Thus $\mathcal{C}=\left(1, n_{R}+1, n_{R}, r-n_{R}, t r, \operatorname{tr}-1\right)$, i.e. it is a basic configuration.

Case 4. $\Phi=-1$ and $\delta=1$.

Then $\mathcal{C}^{\prime}=\left(\tilde{\jmath}-(s-t), 1, s-t, t, n_{C} s, n_{R} s\right)$. By applying Lemma 4.17 to $P B Q$ instead of $A$ we obtain that $n_{C}=n_{R}$. Since for $\mathcal{C}^{\prime}$ the last parameter is $n_{R} s$ and not $n_{R} s-1, \mathcal{C}^{\prime}$ is of the same form of $\mathcal{C}$ as in either case 1 or case 2 (the two cases with $\delta=0$ ). But we excluded these cases already.

## 5. Higher level matrices

In this section, we address the following questions:

- Are there simple composition techniques for constructing high level thin Lehman matrices from low level thin Lehman matrices?
- Are there thin Lehman matrices of arbitrarily high level?
5.1. Compositions. We describe ways of composing Lehman matrices to obtain more complicated, potentially higher level, Lehman matrices.

Proposition 5.1. Let $A, B \in M_{n}(\mathbb{B}), \Sigma_{A}, \Sigma_{A^{\prime}}, \Sigma_{B}, \Sigma_{B^{\prime}} \in M_{n}(\{0, \pm 1\})$ such that $(A, B),(A+$ $\left.\Sigma_{A}, B+\Sigma_{B}\right),\left(A+\Sigma_{A^{\prime}}, B+\Sigma_{B^{\prime}}\right)$ are all Lehman pairs and $A+\Sigma_{A}+\Sigma_{A^{\prime}}, B+\Sigma_{B}+\Sigma_{B^{\prime}} \in M_{n}(\mathbb{B})$. Then $\left(A+\Sigma_{A}+\Sigma_{A^{\prime}}, B+\Sigma_{B}+\Sigma_{B^{\prime}}\right)$ is a Lehman pair iff

$$
\Sigma_{A} \Sigma_{B^{\prime}}^{T}+\Sigma_{A^{\prime}} \Sigma_{B}^{T}=0
$$

Proof. Since $(A, B),\left(A+\Sigma_{A}, B+\Sigma_{B}\right),\left(A+\Sigma_{A^{\prime}}, B+\Sigma_{B^{\prime}}\right)$ are all Lehman pairs, we have

$$
A \Sigma_{B}^{T}+\Sigma_{A}\left(B+\Sigma_{B}\right)^{T}=0 \text { and } A \Sigma_{B^{\prime}}^{T}+\Sigma_{A^{\prime}}\left(B+\Sigma_{B^{\prime}}\right)^{T}=0
$$

Using these two matrix equations and the fact that $A B^{T}=E+I$, we find that

$$
\left(A+\Sigma_{A}+\Sigma_{A^{\prime}}\right)\left(B+\Sigma_{B}+\Sigma_{B^{\prime}}\right)^{T}=(E+I)+\Sigma_{A} \Sigma_{B^{\prime}}^{T}+\Sigma_{A^{\prime}} \Sigma_{B}^{T} .
$$

Therefore, $\left(A+\Sigma_{A}+\Sigma_{A^{\prime}}, B+\Sigma_{B}+\Sigma_{B^{\prime}}\right)$ is a Lehman pair iff $\Sigma_{A} \Sigma_{B^{\prime}}^{T}+\Sigma_{A^{\prime}} \Sigma_{B}^{T}=0$, as desired.

Corollary 5.2. Let $A, B \in M_{n}(\mathbb{B}), \Sigma_{A}, \Sigma_{A^{\prime}}, \Sigma_{B}, \Sigma_{B^{\prime}} \in M_{n}(\{0, \pm 1\})$ such that $(A, B),(A+$ $\left.\Sigma_{A}, B+\Sigma_{B}\right),\left(A+\Sigma_{A^{\prime}}, B+\Sigma_{B^{\prime}}\right)$ are all Lehman pairs and $\operatorname{supp}\left(\Sigma_{A}\right) \bigcap \operatorname{supp}\left(\Sigma_{A^{\prime}}\right)=\emptyset$, $\operatorname{supp}\left(\Sigma_{B}\right) \bigcap \operatorname{supp}\left(\Sigma_{B^{\prime}}\right)=\emptyset$. Then $\left(A+\Sigma_{A}+\Sigma_{A^{\prime}}, B+\Sigma_{B}+\Sigma_{B^{\prime}}\right)$ is a Lehman pair iff

$$
\Sigma_{A} \Sigma_{B^{\prime}}^{T}+\Sigma_{A^{\prime}} \Sigma_{B}^{T}=0
$$

Proof. Since, $\Sigma_{A}$ and $\Sigma_{A^{\prime}}$ have disjoint support, $\left(A+\Sigma_{A}+\Sigma_{A^{\prime}}\right) \in M_{n}(\mathbb{B})$ follows. Similarly, $\Sigma_{B}$ and $\Sigma_{B^{\prime}}$ have disjoint support implies $\left(B+\Sigma_{B}+\Sigma_{B^{\prime}}\right) \in M_{n}(\mathbb{B})$. Now, we can apply Proposition 5.1.
5.2. Long cycles. In some sense, the simplest level-1 update is the one given by a configuration in which all blocks are $1 \times 1$. (See simple- $C 4$ in Section 4.) There is a nice generalization of this simple combinatorial structure to an arbitrary level. We call the general structure $2 \delta$-cycle, for
$\delta \in\{2,3, \ldots, s-1\}$. We define the underlying update by describing the primal perturbation $\Sigma_{A}$ and the dual perturbation $\Sigma_{B}$.

The nonzero entries of $\Sigma_{A}$ are given as follows:

$$
\begin{gathered}
\left(\Sigma_{A}\right)_{11}:=-1 ; \quad\left(\Sigma_{A}\right)_{(\delta-1) r+1,1}:=1 ; \\
\left(\Sigma_{A}\right)_{k r+1,(k+1) r}:=-1, \quad\left(\Sigma_{A}\right)_{(k-1) r+1,(k+1) r}:=1 \forall k \in\{1,2, \ldots, \delta-1\} .
\end{gathered}
$$

All nonzero entries of $\Sigma_{B}$ are in the following 2-by-2 block structure:

|  | $\ell r$ | $\ell r+1$ |
| :---: | :---: | :---: |
| $k r$ | -1 | +1 |
| $k r+1$ | +1 | -1 | for all $1 \leq \ell<k \leq \delta$ such that $(\ell+k)$ is odd.

We denote the above matrices by $\Sigma_{A}(\delta)$ and $\Sigma_{B}(\delta)$.

Proposition 5.3. Let $r \geq 2, s \geq 2$ be arbitrary integers and let $n:=r s-1$. Then for every $\delta \in\{2,3, \ldots, s-1\}, A:=C_{r}+\Sigma_{A}(\delta)$ and $B:=D_{s}+\Sigma_{B}(\delta)$ make a thin Lehman pair.

Proof. It is easy to verify that $A, B \in M_{n}(\mathbb{B})$. To verify that $A B^{T}=E+I$, it suffices to check the matrix equation

$$
C_{r}\left[\Sigma_{B}(\delta)\right]^{T}+\Sigma_{A}(\delta) D_{s}^{T}+\Sigma_{A}(\delta)\left[\Sigma_{B}(\delta)\right]^{T}=0
$$

It is easily seen that (restricted to their nonzero rows and columns),

$$
C_{r}\left[\Sigma_{B}(\delta)\right]^{T}=\begin{array}{r|rrrrrrrrr} 
& 2 r & 2 r+1 & 3 r & 3 r+1 & 4 r & 4 r+1 & \cdots & \delta r & \delta r+1 \\
+1 & -1 & 1 & 0 & 0 & -1 & 1 & \cdots & 0 & 0 \\
r+1 & 1 & -1 & -1 & 1 & 1 & -1 & \cdots & -1 & 1 \\
2 r+1 & 0 & 0 & 1 & -1 & -1 & 1 & \cdots & 1 & -1 \\
3 r+1 & 0 & 0 & 0 & 0 & 1 & -1 & \cdots & -1 & 1 \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\
(\delta-1) r+1 & 0 & 0 & 0 & 0 & 0 & 0 & \cdots & 1 & -1
\end{array}
$$

$$
\begin{aligned}
& \Sigma_{A}(\delta) D_{s}^{T}=\begin{array}{r|rrrrrrrrr|} 
& 2 r & 2 r+1 & 3 r & 3 r+1 & 4 r & 4 r+1 & \cdots & \delta r & \delta r+1 \\
\cline { 2 - 8 } & 1 & 1 & -1 & 1 & -1 & 1 & -1 & \cdots & 1 \\
-1 \\
2 r+1 & -1 & 1 & 0 & 0 & 0 & 0 & \cdots & 0 & 0 \\
3 r+1 & 0 & 0 & -1 & 1 & 0 & 0 & \cdots & 0 & 0 \\
& 0 & 0 & 0 & 0 & -1 & 1 & \cdots & 0 & 0 \\
(\delta-1) r+1 & 0 & \vdots & \vdots & \vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\
& 0 & 0 & 0 & 0 & 0 & \cdots & -1 & 1
\end{array} \\
& \Sigma_{A}(\delta)\left[\Sigma_{B}(\delta)\right]^{T}=\begin{array}{r|rrrrrrrrr|} 
& 3 r & 3 r+1 & 4 r & 4 r+1 & 5 r & 5 r+1 & \cdots & \delta r & \delta r+1 \\
\cline { 2 - 8 } & 1 & -1 & 1 & 0 & 0 & -1 & 1 & \cdots & -1 \\
\hline & 1 & -1 & -1 & 1 & 1 & -1 & \cdots & 1 & -1 \\
2 r+1 & 0 & 0 & 1 & -1 & -1 & 1 & \cdots & -1 & 1 \\
3 r+1 & 0 & 0 & 0 & 0 & 1 & -1 & \cdots & 1 & -1 \\
\vdots & \vdots & \vdots & \vdots & \vdots & \vdots & \vdots & \ddots & \vdots & \vdots \\
& (\delta-1) r+1 & 0 & 0 & 0 & 0 & 0 & 0 & \cdots & 0 \\
0
\end{array}
\end{aligned}
$$

where we illustrated the last two columns and the last rows of the matrices for $\delta$ odd. Therefore, $A B^{T}=E+I$ and $(A, B)$ is a thin Lehman pair.

What is the level of the thin Lehman matrix $C_{r}+\Sigma_{A}(\delta)$ defined above? A likely answer is $\delta-1$ but we could not prove it. It is easy to see that the level of $C_{r}+\Sigma_{A}(\delta)$ is at most $\delta-1$ : Indeed $\Sigma_{A}(\delta)$ has $\delta$ nonzero rows (and columns). When restricted to its support, this matrix is the node-arc incidence matrix of a circuit on $\delta$ nodes. Therefore, $\operatorname{rank}\left(\Sigma_{A}(\delta)\right)=\delta-1$. Hence, the level of $A$ is at most $(\delta-1)$. Note that the highest possible level of $C_{r}+\Sigma_{A}(\delta)$ is $\max \{r, s\}-2$.

Proving lower bounds is much harder. In the next section, we give a lower bounding technique. Note however that the resulting lower bounds are typically not tight.
5.3. Lower bounding the level of thin Lehman matrices. Let $A \in M_{n}(\mathbb{B})$ be $r$-regular for some $r \geq 2$. We define the simple undirected graph $G_{A}:=\left(V\left(G_{A}\right), E\left(G_{A}\right)\right)$ by

- $V\left(G_{A}\right):=\{i: i$ is a row of $A\}$,
- $i j \in E\left(G_{A}\right) \operatorname{iff}\left|\operatorname{row}_{i}(A) \cap \operatorname{row}_{j}(A)\right|=r-1$.

Then, the maximum degree of any node in $G_{A}$ is at most 2 . Thus, $G_{A}$ can be partitioned into vertex-disjoint paths called segments. We denote by segment $(A)$ the number of segments of $G_{A}$. This parameter is invariant under the isomorphisms of $A$.

Remark 5.4. Let $A$ be as above and let $P$ and $Q$ be $n \times n$ permutation matrices. Then

$$
\operatorname{segment}(A)=\operatorname{segment}(P A Q)
$$

Lemma 5.5. Let $A, P$, and $Q$ be as above. Define $\Sigma:=P A Q-C_{r}, t:=\operatorname{rank}(\Sigma)$. Then $\Sigma$ has at most 2 tr non-zero rows.

Proof. Suppose for a contradiction that $\Sigma$ has more than $2 \operatorname{tr}$ non-zero rows. Let $S$ be a minimal set of columns of $\Sigma$ such that the union of their supports covers all non-zero rows of $\Sigma$.

Claim 1. $|S| \geq t+1$.

Proof. By definition, the number of " -1 "s as well as the number of " +1 "s in each column of $\Sigma$ is at most $r$. So,

$$
\left|\operatorname{supp}\left[\operatorname{col}_{j}(\Sigma)\right]\right| \leq 2 r, \text { for all } j .
$$

Thus, $|S| \geq t+1$ as desired.

Claim 2. $\operatorname{col}_{j}(\Sigma)$ for $j \in S$ are linearly independent.

Proof. For every column $j \in S$, the minimality of $S$ implies that there exists a row $i(j)$ that is covered by column $j$ only. Consider the submatrix of $\Sigma$ indexed by the column-row pairs $(j, i(j))$. This submatrix is the $|S| \times|S|$ identity matrix.

We have $\operatorname{rank}(\Sigma) \geq|S| \geq t+1$ (where the first inequality uses Claim 2 and the second uses Claim 1), a contradiction.

Lemma 5.6. Let $\Sigma \in M_{n}(\{0, \pm 1\})$ be 0 -regular with $q$ non-zero rows. Then $C_{r}+\Sigma$ has at most $2 q$ segments.

Proof. Note that $G_{C_{r}}$ is the $n$-circuit. The next elementary observation is all we need.

Claim 1. Let $A \in M_{n}(\mathbb{B})$ be r-regular. Also let $\ell \in\{0, \pm 1\}^{n}$ be 0 -regular. Let $e_{i}$ denote the ith unit vector. Then the only edges in $G_{A}$ possibly not in $G_{A+e_{i} \ell^{T}}$ are incident to vertex $i$.

We apply the above claim repeatedly, starting with $G_{C_{r}}$. There are at most $2 q$ edges of $G_{C_{r}}$ that are not in $G_{C_{r}+\Sigma}$. Since edges of $G_{C_{r}+\Sigma}$ that are not in $G_{C_{r}}$ can only decrease the total number of segments, $G_{C_{r}+\Sigma}$ has at most $2 q$ segments.

Proposition 5.7. Let $A \in M_{n}(\mathbb{B})$ be a thin Lehman matrix that is $r$-regular. Then

$$
\operatorname{level}(A) \geq \frac{\operatorname{segment}(A)}{4 r}
$$

Proof. Let $t:=\operatorname{level}(A)$. Then, there exist $n \times n$ permutation matrices $P, Q$ such that $\Sigma:=$ $P A Q-C_{r} \in M_{n}(\{0, \pm 1\})$ is 0-regular and has rank $t$. Now, Lemma 5.5 implies that $\Sigma$ has at most $2 t r$ non-zero rows. Lemma 5.6 implies that

$$
\operatorname{segment}(P A Q) \leq 4 t r .
$$

Using Remark 5.4 we conclude $t \geq \operatorname{segment}(A) /(4 r)$.

Theorem 5.8. There exist thin Lehman matrices of arbitrarily high level.

Proof. We let $r:=3$ and for large integers $s$, set $n:=r s-1$. We define $A$ from $C_{r}$ by applying the configurations

$$
(1,2,1,1,3,3),(6,7,1,1,3,3),(11,12,1,1,3,3),(16,17,1,1,3,3), \cdots
$$

It is easy to verify that $A$ is a Lehman matrix. Indeed the dual of $A$ is defined from $D_{s}$ by applying the configurations

$$
(2,3,1,1,1,1),(7,8,1,1,1,1),(12,13,1,1,1,1),(17,18,1,1,1,1), \cdots
$$

as can be checked by multiplying these two matrices. Consider those integers $n$ satisfying the above condition and $n=5 k$, for some integer $k \geq 4$. Then segment $(A) \geq 2 k$. Using Proposition 5.7, we conclude that

$$
\operatorname{level}(A) \geq \frac{n}{30}
$$

Therefore, $\operatorname{level}(A)=\Omega(n)$ for this construction.

Remark 5.9. Consider the long cycle construction. Let $A$ be as defined in Proposition 5.3. It is easy to check that a $2 \delta$-cycle creates $\delta$ segments, the largest value $\delta$ can take is $s-1$. Thus, Proposition 5.7 implies

$$
\operatorname{level}(A) \geq \frac{s-1}{4 r}
$$

for the largest value of $\delta$. If $r=3$, then $3 s=n+1$ and the long cycle construction also yields a proof of Theorem 5.8:

$$
\operatorname{level}(A) \geq \frac{n-2}{36}
$$

## 6. Fat matrices

### 6.1. Examples.



Matrices $F_{7}$ and $P_{10}$ are fat Lehman matrices. Matrix $F_{7}$ is the point-line incidence matrix of the Fano plane. $F_{7}$ is self-dual, thus $k=2$ in (1). Matrix $P_{10}$ is the matrix whose columns correspond to the edges of $K_{5}$ and whose rows are the incidence vectors of the triangles of $K_{5}$. Equivalently, $P_{10}$ can be viewed as the vertex-vertex incidence matrix of the Petersen graph (hence the notation). $P_{10}, P_{10}+I$ form a Lehman pair, thus $k=2$ in (1).

### 6.2. Determinant.

Remark 6.1. In this section $E_{n}$ denotes the $n \times n$ matrix of 1 s. For $n \geq 2$ and $k \geq 1$, the matrix $E_{n}+k I_{n}$ has two distinct eigenvalues, namely $k$ with multiplicity $n-1$, and $n+k$ with multiplicity 1 . In particular,

$$
\operatorname{det}\left(E_{n}+k I_{n}\right)=k^{n-1}(n+k)
$$

Proof. Since $\left(E_{n}+k I_{n}\right)-k I_{n}=E_{n}$ and there are $n-1$ linearly independent vectors in $\operatorname{Null}\left\{e_{n}\right\}$, the multiplicity of $k$ is at least $n-1$. Vector $e_{n}$ is the eigenvector for the eigenvalue $n+k$. Since the total multiplicity is at most $n$, the result about eigenvalues follows. Finally, the determinant is the product of the eigenvalues.

As an example, consider $F_{7}$ in (5). Then $n=7$ and since $F_{7}$ is self-dual, $k=2$. Hence $\operatorname{det}\left(E_{7}+2 I_{7}\right)=9 \times 2^{6}$ and $\operatorname{det}\left(F_{7}\right)=3 \times 2^{3}$.

Remark 6.2. Let $A$ be an $r$-regular Lehman matrix.
(i) If $A$ is thin, then $|\operatorname{det}(A)|=r$,
(ii) If $A$ is self-dual, then $|\operatorname{det}(A)|=(r-1)^{\frac{r(r-1)}{2}} r$.

Proof. (i) By Theorem 2.1, the dual of $A$ is an $s$-regular matrix $B$ such that $r s=n+1$. Remark 6.1 implies that $\operatorname{det}\left(E_{n}+I_{n}\right)=n+1=r s$. Thus $\operatorname{det}(A) \operatorname{det}(B)=\operatorname{det}\left(E_{n}+I_{n}\right)=r s$.

Since $A$ is an $r$-regular nonsingular integral matrix, it follows that its determinant is a nonzero integer multiple of $r$. Thus $|\operatorname{det}(A)| \geq r$ and similarly $|\operatorname{det}(B)| \geq s$, and the result follows.
(ii) Since $A$ is self-dual, $k=r-1$. By Theorem 2.1, $r^{2}=n+r-1$. Remark 6.1 implies that $\operatorname{det}(A)^{2}=\operatorname{det}\left(E_{n}+(r-1) I_{n}\right)=(r-1)^{r(r-1)} r^{2}$. The result follows.

Recall that $|\operatorname{det}(A)|$ equals the volume of the parallelopiped defined by the columns of $A$ (viewed as vectors of $\mathbb{R}^{n}$ ). This justifies our terminology of thin Lehman matrix (the parallelopiped formed by its columns has the smallest possible volume among all nonsingular $r$-regular matrices in $M_{n}(\mathbb{B})$ ). By contrast, fat Lehman matrices give rise to parallelopipeds with larger volumes, the extreme case being that of nondegenerate finite projective planes.
6.3. Lehman matrices from projective planes. A projective plane consists of points and lines such that any two distinct points belong to exactly one line, and any two distinct lines intersect in exactly one point. A projective plane is degenerate if at least three of any four points belong to the same line. It can be shown that all the lines of a nondegenerate finite projective plane have the same number of points. Therefore, point-line incidence matrices $A \in M_{n}(\mathbb{B})$ of nondegenerate finite projective planes are exactly the solutions of the equation $A A^{T}=E+k I$, i.e. they are the self-dual Lehman matrices. We review known results about these matrices. First note that Theorem 2.1 implies that $n=k^{2}+k+1$. The integer $k$ is called the order of the projective plane.

Not all orders $k$ are possible, as proved by Bruck and Ryser [2] in the following theorem.

Theorem 6.3. If $k=1,2(\bmod 4)$ and $y^{2}+z^{2}=k$ has no solution in integers, then there is no projective plane of order $k$.

For example, this implies that there are no projective planes of orders 6 and 14 . What is the idea of the proof of the Bruck-Ryser theorem? Observe that $E+k I$ is a positive definite matrix. Therefore it always has a decomposition $A A^{T}=E+k I$. Bruck and Ryser [2] address the question of whether there exists such a decomposition where $A$ has rational entries. (When
$n=1$, this question reduces to: When does there exist a rational number $a$ such that $a^{2}=1+k$ ?) By clever arguments, Bruck and Ryser massage the quadratic form $x^{T} A A^{T} x=x^{T}(E+k I) x$ (which has nonzero rational solutions) until they eventually reduce it to $y^{2}+z^{2}=k$ in integers.

Does this line of proof carry over to the general Lehman equation $A B^{T}=E+k I$, i.e. can we use the fact that $A$ and $B$ have rational entries to exclude certain values of $k$ ? Unfortunately not: For any nonsingular rational matrix $A$, we can set $B^{T}=A^{-1}(E+k I)$ which is also rational. In order to prove the nonexistence of Lehman matrices for certain values of $k$, one needs combinatorial arguments using the fact that $A, B$ are 0,1 matrices.

The following table gives the number of projective planes for small orders $k$.

| $k$ | 2 | 3 | 4 | 5 | 6 | 7 | 8 | 9 | 10 | 11 | 12 | 13 | 14 | 15 | 16 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| Number | 1 | 1 | 1 | 1 | 0 | 1 | 1 | 4 | 0 | $\geq 1$ | $?$ | $\geq 1$ | 0 | $?$ | $\geq 22$ |

Next we describe an infinite family of projective planes denoted by $\operatorname{PG}(2, k)$. Let $V$ be a 3-dimensional vector space over a finite field with $k$ elements. The points of $\operatorname{PG}(2, k)$ are the 1-dimensional subspaces of $V$ and its lines are the 2-dimensional subspaces of $V$. Then $\mathrm{PG}(2, k)$ is a projective plane of order $k$. For example, when $k=2$ we get the Fano plane $F_{7}$.

This construction implies that a projective plane of order $k$ exists whenever $k$ is a prime power, since there always exists a finite field with $k$ elements in this case. Interestingly, all known examples of finite projective planes have an order which is a prime power.
6.4. Nearly self-dual Lehman matrices. We call nearly self-dual a Lehman matrix $A$ with the following properties:
(i) $A=A^{T}$ and
(ii) the dual of $A$ is $A+I$.

Theorem 6.4. Let A be a nearly self-dual Lehman matrix which is $r$-regular. Then $r=2,3,7$ or 57 .

Proof. Since $A+I$ is a 0,1 matrix, the entries in the diagonal of $A$ are all equal to 0 . Since $A=A^{T}$, the matrix $A$ is the vertex-vertex incidence matrix of a graph $G$. Since $A$ is $r$-regular, $G$ is $r$-regular.

## Claim 1. The graph $G$ has girth at least 5.

Proof. Suppose otherwise. Then $G$ contains a triangle with vertices $i, j, k$ or a 4 -cycle with vertices $i, k, j, l$ in that order. In both cases, the scalar product $\left\langle\operatorname{row}_{i}(A), \operatorname{row}_{j}(A+I)\right\rangle \geq 2$. But this contradicts Lehman's equation, which implies $\left\langle\operatorname{row}_{i}(A), \operatorname{row}_{j}(A+I)\right\rangle=1$ for $i \neq j$.

An $(r, g)$-cage is a graph that (i) is $r$-regular, (ii) has girth at least $g$, and has the smallest possible number of vertices among all graphs satisfying (i) and (ii).

Claim 2. The graph $G$ is an $(r, 5)$-cage with $1+r^{2}$ vertices.

Proof. Consider any $r$-regular graph $H$ with girth at least 5, and let $v$ be a vertex of $H$. Vertex $v$ has $r$ neighbors $v_{1}, \ldots, v_{r}$ and each of these vertices $v_{i}$ has $r-1$ neighbors distinct from $v$. Furthermore, all these vertices are distinct since $H$ contains no 4-cycle. Therefore, $H$ has at least $1+r+r(r-1)=1+r^{2}$ vertices.

Since $A, A_{I}$ is a Lehman pair, it follows from Theorem 2.1 (ii) that $r(r+1)=n+(r-1)$, i.e. the graph $G$ has $n=1+r^{2}$ vertices. Thus $G$ is an $(r, 5)$-cage.

A theorem of Hoffman and Singleton [6] states that, for any ( $r, 5$ )-cage, $n \geq 1+r^{2}$ and equality holds if and only if $r \in\{2,3,7,57\}$.

Hoffman and Singleton [6] show that there is a unique solution (up to isomorphism) for each of the cases $r=2,3,7$. The existence of a solution for the case $r=57$ is unknown.

The case $r=2$ (i.e. $n=5$ ) is the circulant $C_{2}^{5}$.
The case $r=3$ (i.e. $n=10$ ) is the Petersen matrix $P_{10}$ mentioned earlier.
The case $r=7$ (i.e. $n=50$ ) was constructed by Hoffman and Singleton [6].
6.5. Fat Lehman matrices and minimally nonideal matrices. The point-line matrices of degenerate finite projective planes are minimally nonideal. The cores of most other known minimally nonideal matrices are thin Lehman matrices. We know only three exceptions: $F_{7}, P_{10}$ and its dual. These three fat Lehman matrices play a central role in Seymour's conjecture about ideal binary matrices [13]. A 0,1 matrix is binary if the sum modulo 2 of any three of its rows is greater than or equal to at least one row of the matrix. Seymour's conjecture states that there are only three minimally nonideal binary matrices $\left(F_{7}, \mathcal{O}_{K_{5}}\right.$ whose columns are indexed by the edges of $K_{5}$ and whose rows are the characteristic vectors of the odd cycles of $K_{5}$, and its blocker): Their cores are $F_{7}, P_{10}$ and its dual respectively.

## 7. OPEN PROBLEMS AND CONCLUDING REMARKS

The Lehman matrix equation (1) occurs prominently in the study of minimally nonideal matrices. Bridges and Ryser [1] give basic properties of its solutions (Theorem 2.1). Two infinite families of solutions are known: thin Lehman matrices and finite projective planes. In this paper, we classify thin Lehman matrices according to their similarity to the circulant matrices $C_{r}^{n}$ : Level $t$ matrices are isomorphic to $C_{r}^{n}$ plus a rank $t$ matrix. We were able to describe explicitly all level 1 matrices and we showed that level $t$ matrices can be described by a number of parameters that only depends on $t$ (independent of $n$ and $r$ ). We also gathered results from the literature that are relevant to our understanding of fat Lehman matrices. There remain many open problems.

Question 1: Are there other infinite families of Lehman matrices beside thin matrices and projective planes?

Question 2: Can Theorem 1.2 be strengthened as follows: If $A$ is a thin Lehman matrix of level $t$, then $A$ can be described with $O(t)$ parameters?

In particular, can every thin $n \times n$ matrix be described with only $O(n)$ parameters?

Question 3: Do all thin Lehman matrices have level at most $\frac{n}{\min (r, s)}$ ?

Question 4: Is there a decomposition theorem stating that a thin Lehman matrix either is in a well-described family (such as matrices with low level or long cycles) or has a decomposition (such as presented in Section 5)?

Question 5: Is a thin Lehman matrix always the core of some minimally nonideal matrix?

Question 6: Is $F_{7}$ the only nondegenerate finite projective plane whose point-line matrix is the core of a minimally nonideal matrix? Beth Novick [12] answered this question positively when "the core of" is removed from the statement.

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