Davenport–Schinzel theory of matrices

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Abstract


Let C be a configuration of 1’s. We define \( f(n; C) \) to be the maximal number of 1’s in a 0–1 matrix of size \( n \times n \) not having \( C \) as a subconfiguration. We consider the problem of determining the order of \( f(n; C) \) for several forbidden \( C \)'s. Among other results we prove that \( f(n; (1, 1)) = \Theta(a(n)n) \), where \( a(n) \) is the inverse of the Ackermann function.

1. Introduction

A configuration, \( C = (c_{ij}) \) \( (1 \leq i \leq u, 1 \leq j \leq v) \), is a partial matrix with 1’s and blanks at the entries. All the matrices we are going to work with will be 0–1 matrices. We say that a matrix \( M = (m_{ij}) \) does have the configuration \( C \) if one can find \( u \) rows \( i_1, i_2, \ldots, i_u \), \( i_1 < \cdots < i_u \) and \( v \) columns \( j_1, j_2, \ldots, j_v \), \( j_1 < \cdots < j_v \) in \( M \) such that the corresponding submatrix contains \( C \), i.e., \( m_{i_k,j_{\ell}} = 1 \) whenever \( c_{i_k,j_{\ell}} = 1 \). Let \( f(n, m; C) \) denote the maximum number of 1’s in an \( n \times m \) matrix \( M \) not containing \( C \). In the case of \( n = m \) we write \( f(n; C) \). One can allow several forbidden configurations, the corresponding threshold function is \( f(n, m; \{C_1, \ldots, C^n\}) \) or \( f(n; \{C_1, \ldots, C^n\}) \).

Our research is closely related to previous works in combinatorics.

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First, let us mention Turán’s theory in extremal graph theory. There the question is: Given a graph $G$, what is $T(n; G)$, the maximum number of edges of a graph with $n$ vertices and not containing $G$ as a subgraph? A special case is when we work in the universe of bipartite graphs. Our matrices can be considered as bipartite graphs. The important difference between Turán’s theory and our question that in our case the vertices (the rows and columns) are ordered. This is a very important difference but in some special case the restriction on the order is insignificant. An example is the four cycle (complete bipartite graph between two color classes of size 2 each). Classical results in graph theory [14, 8, 5] immediately give us the following theorem.

**Theorem 1.1.** $f\left(n; \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}\right) = \Theta(n^{3/2})$.

We do not know exactly how these two problems are related, but the following facts are known: The Erdős–Stone–Simonovits Theorem ([9, 10], for a survey see Bollobás’ book [4]) says that the order of magnitude of $T(n; G)$ depends on the chromatic number of $G$, namely

$$\lim_{n \to \infty} \frac{T(n; G)}{\binom{n}{2}} = 1 - (\chi(G) - 1)^{-1}.$$  

This theorem gives a sharp estimate on $T(n; G)$, except for bipartite $G$. For every bipartite graph $B$ which is not a tree there are positive constants $c_1$ and $c_2$ (not depending on $n$) such that

$$\Omega(n^{1+c}) \leq T(n; B) \leq O(n^{2-c})$$

holds. If the graph is a tree $F$, then it is straightforward that $T(n; F) = \Theta(n)$. However we will see that our problem has completely different threshold functions. For a special matrix (such that the corresponding graph is a tree, hence it has linear Turán function) our threshold function turns out to be $\Theta(n \log n)$.

An other related question is raised by Davenport and Schinzel. A sequence $s = x_1x_2 \cdots x_t$ is called a *Davenport–Schinzel sequence*, $s \in DS_k(n)$, if $x_i \neq x_{i+1}$, $x_i \in \{1, 2, \ldots, n\}$ and $s$ does not contain a subsequence $x_{i_1}x_{i_2} \cdots x_{i_k}$ such that

$$x_{i_1} = x_{i_3} = \cdots = x_{i_{2i-1}} = \cdots \neq x_{i_2} = x_{i_4} = \cdots = x_{i_{2i}} = \cdots$$

$(i_1 < i_2 < \cdots < i_k)$. Let $ds_k(n)$ denote the maximum length of an $s \in DS_k(n)$. It is obvious that

$$ds_1(n) = n, \quad ds_2(n) = 2n - 1.$$  

Szemerédi [17] proved that $ds_k(n) = O(n \log^* (n))$ for all fixed $k$ while $n$ tends to infinity. (Here, as usual, $\log^* n$ denotes the inverse of the function $p : \mathbb{N} \to \mathbb{N}$ with $p(1) = 2$, $p(n + 1) = 2^{p(n)}$.) Recently, mainly due to the works of Sharir [16, 12, 13] it is known that the true order of the magnitude of $ds_k(n)$ for $k \geq 5$ is
really superlinear, e.g., [12]
\[ ds_2(n) = \Theta(n\alpha(n)), \]
where \( \alpha(n) \) is the inverse Ackermann function, a very slowly growing function. For more on this see Section 7 and 8.

Some specific configurations were investigated in previous papers (see [6, 11]). The motivation of those results were geometrical.

For a matrix \( M \) (or vector as a special case) \(|M|\) denotes the number of its entries equal to 1, \( M^T \) is its transpose. \([n]\) is the set of the first \( n \) positive integers, and \([a, b] =: \{a, a+1, \ldots, b\} \).

2. A reduction between matrices

Let \( C \) be a configuration of 1's. We are going to define two operations on \( C \). The first one is simply deleting an entry. The second one is attaching a new column or row to the boundary of \( C \) and placing an entry 1 in the new column or row, next to an existing one in \( C \).

**Definition 2.1.** If \( D \) can be constructed from \( C \) using one of these operations we say that \( D \) is obtained by an *elementary operation* from \( C \). We use the notation \( C \rightarrow^\ast D \). Let \( \rightarrow \) be the transitive closure of \( \rightarrow^\ast \), i.e., \( C \rightarrow D \) if \( D \) can be constructed from \( C \) using a sequence of elementary operations.

Note that the size of the matrix can decrease by the first type of elementary operation if the deletion of the given entry creates an empty row or column. Fig. 1 shows several configurations and their relations.

**Theorem 2.2.** Let \( C, D \) be configurations such that \( C \rightarrow D \) by \( t \) elementary steps. Then \( f(n, m; D) \leq f(n, m; C) + t \cdot \max(n, m) \).

**Proof.** It is sufficient to prove the case \( t = 1 \), we can assume that \( C \rightarrow^\ast D \). If \( D \) is constructed by deleting an entry then the claim is obvious. So we can assume that \( D \) is constructed by adding an extra column to the end of \( C \) with an extra 1 (the other cases are very similar). Let \( M \) be a matrix of size \( n \times m \) with \( f(n, m; D) \) many 1's such that it does not have \( D \) as a subconfiguration. Let \( M' \) be the matrix that we get if we delete the last 1 in each row (assuming that there is any). Easy to realize that \( M' \) does not have \( C \) as a subconfiguration. So the number of remainder 1's in \( M' \) is at most \( f(n, m; C) \). \( \square \)

The natural way to apply Theorem 2.2 is that in the case of \( C \rightarrow D \), an upper bound on \( f(n; C) \) gives an upper bound on \( f(n; D) \) and a construction for a matrix not having \( D \) as a submatrix gives a good construction for \( C \).

Fig. 2 contains some additional matrices with four 1's and some of their \( \rightarrow \) relations.

Let \( B_2 \) be \((1, 1)\), a 1 \( \times \) 2 configuration.
\[ C_1 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_2 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_3 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_4 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_5 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_6 = \begin{pmatrix} 1 & 1 & 1 \\ 1 & 1 & 1 \end{pmatrix} \]

\[ C_7 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_8 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_9 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_{10} = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

**Proposition 2.3.** If \( D_2 \rightarrow C \) and \( C \) has at least 2 entries in it then
\[
\min(n, m) \leq f(n, m; C) \leq c_C(n + m).
\]

**Proof.** Trivial. The lower bound comes considering a matrix \( M \) with 1's only in one row or in one column.

The upper bound is immediate from Theorem 2.2. \( \square \)

\[ C_{11} = \begin{pmatrix} 1 & 1 & 1 \\ 1 & 1 & 1 \end{pmatrix} \]

\[ C_{12} = \begin{pmatrix} 1 & 1 \\ 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_{13} = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_{14} = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \]

\[ C_{15} = \begin{pmatrix} 1 \\ 1 \end{pmatrix} \]

Fig. 1.

Fig. 2.
Fig. 3.

We remark that Figs. 1–3 contain all the 37 configurations with four 1’s (not distinguishing two if they are the same up to rotations and reflections). The simple reduction principle yields that 22 of them have linear complexity.

**Corollary 2.4.** (1) If $M$ has at most 3 nonzero entries then $f(n, m; M) \leq 2(n + m)$.

(2) The 22 matrices in Fig. 3 have linear complexity, $f(n, m; C_i) \leq 3(n + m)$ for $16 \leq i \leq 37$.

One can extend the $\rightarrow$ relations to sets of configurations. This will be proven very useful.

**Definition 2.5.** Let $C^1, \ldots, C^k$ be a set of configurations. We are going to define two operations. One is simply adding a new configuration to our set. The second is substitute a $C'$ with $D$ if $C' \rightarrow D$. The transitive closure of these relations is $\rightarrow$.

The notation is not in conflict with Definition 2.1, which is a special case of this. Note that $\{C^1, \ldots, C^k\} \rightarrow \{D^1, \ldots, D^l\}$ iff for every $i$ there is a $j$ such that $C^i \rightarrow D^j$ according to Definition 2.1.

The analog of Theorem 2.2 is the following.
Theorem 2.6. If \((C^1, \ldots, C^k) \rightarrow (D^1, \ldots, D^l)\) then
\[
f(n, m; (D^1, \ldots, D^l)) \leq f(n, m; (C^1, \ldots, C^k)) + \text{const}(n + m),
\]
where the constant depends only on the two systems, and not on \(n\) and \(m\).

A few examples:

\[
\begin{pmatrix}
1 & 1 \\
1 & 1
\end{pmatrix}, \begin{pmatrix}
1 & 1 \\
1 & 1
\end{pmatrix}, \begin{pmatrix}
1 & 1 \\
1 & 1
\end{pmatrix} \rightarrow \begin{pmatrix}
1 & 1 \\
1 & 1
\end{pmatrix},
\]
\[
\begin{pmatrix}
1 & 1 \\
1 & 1
\end{pmatrix}, \begin{pmatrix}
1 & 1 \\
1 & 1
\end{pmatrix} \rightarrow \begin{pmatrix}
1 & 1 \\
1 & 1
\end{pmatrix}.
\]

3. Matrices with \(n \log n\) complexity

Theorem 3.1 [11]. \(f(n, \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}) < 6n \log n\).

The construction in [11] shows that this upper bound is the best up to a constant factor. Below we give another, a simpler recursive construction.

Construction 3.2. Let
\[
A_1 = \begin{pmatrix} 1 & 1 \\ 1 & 0 \end{pmatrix}, \quad A_{n+1} = \begin{pmatrix} E_{2^n} & A_n \\ A_n & 0_{2^n} \end{pmatrix},
\]
where \(E_n\) is an \(n \times n\) matrix with 1's only in the diagonal connecting the upper right to the lower left corner, and \(0_n\) the \(n \times n\) zero matrix.

Claim 3.3. (1) \(A_n\) is a \(2^n \times 2^n\) matrix with \((n + 2)2^{n-1}\) many 1's.

(2) \(A_n\) does not have
\[
C_4 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}
\]
as a subconfiguration.

Proof. (1) Easy induction.

(2) Using induction. The initial case is obvious. Let us assume that the claim is verified for \(A_k\), when \(k < n\).
Suppose on the contrary that $A_n$ has the forbidden configuration. $A_n$ is, by
definition, divided into 4 submatrices. We distinguish different cases depending
on which submatrix has the upper left corner of the forbidden configuration. If
one of the $A_{n-1}$'s is the one, then our inducional hypothesis gives the
contradiction. If $E_{2n}$ has that entry then it is easy to verify that the bottom right
corner of the configuration must be in $O_{2n}$. This contradicts the fact that $O_{2n}$ has
no 1 entry at all. □

**Corollary 3.4.** (1) $f(n; C_2), f(n; \{C_2, C_2^T\}), f(n; C_4) = \Theta(n \log n)$.
(2) $f(n; C_i) < 10n \log n$, for $4 \leq i \leq 15$.

**Proof.** (1) Both the lower and upper bound comes from the following relations:

$$C_2 \rightarrow \{C_2, C_2^T\} \rightarrow C_4.$$

(2) See Figs. 1 and 2. □

### 4. A construction with $\left(\frac{n \log n}{\log \log n}\right)$ 1's

In the previous section we saw an $n \log n$ upper bound on $f(C_5)$. Now we
construct a matrix with $\Theta((n \log n)/(\log \log n))$ 1's and not having $C_5$ as a
subconfiguration. This section is a slightly simplified version of [6]. Our
construction will be recursive and it defines $N(s, t)$, a matrix of size $st \times st$, where
$s, t \geq 1$.

First we discuss a few properties of $N(s, t)$ which we need for the formal
definition of the matrix. The $st$ rows are divided into $s$ blocks, each having $t$
consecutive rows. In each block we have a column such that each of its entries are
1's and these are the first 1's in the corresponding rows. This column is the
**leading column** of that block.

Let $N(s)$ be an $s \times s$ matrix without the configurations:

$$C_5 = \begin{pmatrix}
1 & 1 \\
1 & 1 \\
1 & 1 \\
\end{pmatrix}, \begin{pmatrix}
1 & 1 \\
1 & 1 \\
\end{pmatrix}, \begin{pmatrix}
1 & 1 \\
1 & 1 \\
\end{pmatrix}, \begin{pmatrix}
1 & 1 \\
1 & 1 \\
\end{pmatrix}. \tag{4.1}
$$

**Definition 4.1.** $N(1, t)$ is a $t \times t$ matrix with $t$ 1's in the first column and 0's
everywhere else. $N(s, 1)$ is the $s \times s$ identity matrix.

The construction of $N(s, t + 1)$ is the following (we assume that $N(s, t')$ and
$N(s', t')$ are already constructed for $t' \leq t$, $s' < s$ and $t''$ arbitrary): Take a copy of
$N(s, t)$ and insert an **extra row** after each block. In each extra row put a 1 at the
leading column of the block just above it. Add **new columns** at the end of the
already constructed part. At the intersection of extra rows and new columns we
have an $s \times s$ space. Put a copy there of $N(s)$ with maximum number of ones.

The promised properties are maintained so our recursion is correct.
Theorem 4.2. $N(s, t)$ does not have the configurations given in (4.1).

Proof. An easy induction by case by case check. □

The previous theorem gives lower bounds on the complexity of several configurations and sets of configurations.

Corollary 4.3 [6]. $f(n; C_s) = \Omega(n \log n / \log \log n)$.

Proof. Let $f(s, t) = \|N(s, t)\|$ and $f(s) = \max \|N(s)\|$. We have

$$f(s, t + 1) \geq f(s, t) + f(s) + s,$$

and for $s \geq ab$, $f(s) \geq f(a, b)$. These inequalities imply that

$$f(l^a, t) \geq (t - 1)(l^a + l^{a-1} - (l-1)^a) + l^a,$$

especially

$$f(l^{a+1}) \geq f(l^a, l) \geq l^{a+2} - l(l-1)^{a+1}.$$

Letting $a = l - 1$, $n = l'$ we obtain the desired bound. □

5. More matrices with linear complexity

Recall that

$$C_{11} = \begin{pmatrix} 1 & 1 \\ 1 & 1 \\ 1 & 1 \end{pmatrix}.$$

In this section we prove, that the complexities of $C_{11}, \ldots, C_{15}$ are all linear, at most $9n$. As one can see from Fig. 2, and Theorem 2.3 the above result is implied by the following theorem.

Theorem 5.1. $f(n, C_{11}) \leq 7n$.

Proof. Let $A' = (a_{ij})$ be an $n \times m$ 0–1 matrix without $C_{11}$. Delete the first and the last entry in each row, and delete all entries in that row if $\|a'_{ij}\|_{1 \leq j \leq m} \leq 3$. For the obtained matrix $A = (a_{ij})$ we have

$$\|A'\| \leq \|A\| + 3n. \quad (5.2)$$

$A$ does not contain the following configurations either:

$$\begin{pmatrix} 1 & 1 \\ 1 & 1 \\ 1 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 1 \end{pmatrix}.$$
For the $i$th row (if it is non-empty) let $m(i)$ ($M(i)$) denote the minimum (maximum, resp.) index in that row, i.e., $m(i) := \min\{j: a_{ij} = 1\}$. Then $[m(i), M(i)] \subset [m(i'), M(i')]$ implies $i \leq i'$.

The element $a_{ij}$ is called type $\alpha$ if $a_{ij} = 1$, it is not the first neither the last one in its row ($m(i) < j < M(i)$), $\alpha < i$, $j \in [m(\alpha), M(\alpha)]$, and $i$ is minimal with respect to these constraints. By definition, there are no two entries of type $\alpha$ in distinct rows. But there are no two 1's of type $\alpha$ in the $i$th row neither, otherwise together with $a_{\alpha, m(\alpha)}$ and $a_{\alpha, M(\alpha)}$ they form a forbidden subconfiguration. So the number of entries in $A$ which are:

1. first or last in their row,
2. on the top of their column, or
3. have a type $\alpha$,

is at most $4n$. We claim that all the entries of $A$ fall in one of the above 3 categories, implying $\|A\| \leq 3n + m$. Then (5.2) finishes the proof of the theorem.

Suppose that the entry $a_{ij} = 1$ is not the first or the last one in the $i$th row, and that there exists a $t \leq i$ with $a_{eq} = 1$. Then $j \in [m(t), M(t)]$. Let $\alpha$ be the maximum index, such that $\alpha < j$, and $j \in [m(\alpha), M(\alpha)]$. Then $a_{ij}$ has type $\alpha$.

Indeed, suppose on the contrary, that some entry $a_{i'j'}$ has type $\alpha$, with $\alpha < i' < i$. Then, $j \in [m(\alpha), M(\alpha)] \subset [m(i'), M(i')]$, so the existence of $i'$ contradicts the definition of $\alpha$. \qed

Let $C'$ be a $2 \times (t + 2)$ configuration with 1's in the positions $(1, 1)$, $(1, t + 2)$ and $(2, 2), \ldots, (2, t + 1)$. $C'_{11} = C^2$. Deleting from every row the middle $t - 2$ entries, Theorem 5.1 implies the following.

**Corollary 5.3.** $f(n; \begin{pmatrix} 1 \cdot \cdot \cdot 1 \\ 1 1 \cdot \cdot \cdot 1 \end{pmatrix}) = f(n; C') \leq (t + 5)n$.

Finally we mention a generalization of this idea in the direction of sequences with forbidden subsequences. The following corollary is a special case of the result in [2].

**Corollary 5.4.** Suppose that the sequence $s = x_1, x_2, \ldots, x_l$ with $x_i \in [n]$, does not have two identical consecutive members, and does not contain the subsequence abba, where $a < b$, then $l \leq 100n$.

**Proof (sketch).** Split $s$ into $n$ equal parts $s = s_1, s_2, \ldots, s_n$. $\|s_i\| = 100$. Then there is a subset $s'_i \subset s_i$ containing only distinct elements with $|s'_i| = 9$. Put 1's into the $i$th column of an $n \times n$ matrix $A$ according to $s'_i$. Finally, apply Theorem 5.1 to $A$ to get a $C_{11}$, and then to get an abba in $s$. \qed
6. A covering lemma

In this section we prove a covering lemma about 0–1 matrices. As an easy application of our lemma we get several new matrices with linear complexity.

We start with a definition. An intersection of \( s \) consecutive rows and \( t \) consecutive columns is called a rectangle. The horizontal size of \( R \) is \( t \) and it is denoted by \( h(R) \), the vertical size of \( R \) is \( s \) and it is denoted by \( v(R) \). \( M \) itself, is an example for a rectangle.

**Lemma 6.1.** Let \( M \) be an arbitrary 0–1 matrix. Then there is a system of rectangles \( \{ R_i \} \) such that:

1. the \( R_i \) cover all the 1’s,
2. \( \sum_i h(R_i) \leq 4h(M) \) and \( \sum_i v(R_i) \leq 4v(M) \),
3. each \( R_i \) has a 1 in the upper left or bottom right corner.

**Proof.** Let us define a partial order between the positions in a given matrix. We say that \( a \preceq b \), if the row of \( a \) is not later than \( b \)’s row and \( a \)’s column is not later than \( b \)’s column. \( a \preceq b \) if \( a \preceq b \) and \( a \neq b \).

There are incomparable positions. For two incomparable positions \( c \) and \( d \) we say that \( c \prec d \) if \( c \)’s row is earlier than \( d \)’s.

Take \( M \) and consider only the positions where we have a 1. Let \( m_1 \preceq m_2 \preceq \cdots \preceq m_k \) be the set of minimal 1’s for the partial order \( \preceq \). Let \( M_1 \preceq M_2 \preceq \cdots \preceq M_l \) be the set of maximal 1’s for the partial order \( \preceq \). We can assume that \( m_1 \) is in the first column, \( m_k \) is in the first row, \( M_1 \) is in the last row and \( M_l \) is in the last column of \( M \).

Let \( m_{i+1/2} \) (for \( i = 1, \ldots, k-1 \)) be the position in the intersection of the row of \( m_i \) and the column of \( m_{i+1} \). Let \( m_{i+1/2} \) be the lower left corner of \( M \). Let \( m_{k+1/2} \) be the upper right corner of \( M \). Let \( M_{i+1/2} \) (for \( j = 1, \ldots, l-1 \)) be the position in the intersection of the column of \( M_i \) and the row of \( M_{i+1} \). Let \( M_{i+1/2} = m_{i+1/2} \) and \( M_{i+1} = m_{i+1/2} \). Let \( h_i = [m_i, m_{i+1/2}] \) be a horizontal interval of positions in the row of \( m_i \), with endpoints at \( m_i \) and \( m_{i+1/2} \). Let \( v_i \) be the vertical interval \([m_{i-1/2}, m_i]\). We define the corresponding intervals for maximal 1’s. Let \( V_i = [M_i, M_{i+1/2}] \) and \( H_i = [M_{i-1/2}, M_i] \). It is clear that \( v_1, h_1, v_2, \ldots, v_k, h_k \) and \( H_1, V_1, H_2, V_2, \ldots, H_i, V_i \) define two stair shaped curves. Let us denote them by \( S \) and \( s \). By definition it is straightforward that there is no 1 above \( s \) and below \( S \).

Now we are starting to construct our covering system of rectangles. This system is containing two sequences of rectangles: \( \{ Q_i \} \) and \( \{ P_i \} \). The \( Q_i \)’s are going to have an entry 1 at the bottom left corner, the \( P_i \)’s are going to have a 1 at the upper left corner. We define them recursively.

Let \( Q_1 \) be a rectangle with lower right corner at \( M_1 \), with lower left corner at \( m_{1/2} \). So its right vertical side is on the vertical half line starting at \( M_1 \), going up. The missing corner of \( Q_1 \) on this line is where it first hits \( s \).

\( Q_1 \) might cover several \( h_i \) intervals. Let \( h_i \) be the first one which is not covered
by $Q_1$. Let $P_i$ be a rectangle with upper left corner at $m_i$. This fact gives us two half lines starting at $m_i$ and going down and to the right. They hit $S$ at two positions. They will be two other corners of $Q_1$.

Next, we will explain the general step in the definition.

Let us assume that we already defined $Q_1$, $P_1$, $Q_2$, $P_2$, $Q_3$, $P_3$. Let $V_i$ be the first vertical interval of $S$ which is not covered by $Q_1 \cup \cdots \cup P_i$. Let $M_{ij}$ be the bottom right corner of $Q_{i+1}$. This defines two half lines starting at $M_{ij}$, one going up (let us say $e_{i+1}$) and one going to the left. They hit $S$ at two positions. They give us two other corner of $Q_{i+1}$. This completes the definition of $Q_{i+1}$.

Let us assume that we already defined $Q_1$, $P_1$, $Q_2$, $P_2$, $Q_3$, $P_3$, $Q_{i+1}$, $P_{i+1}$. Let $h_i$ be the first horizontal interval of $S$ which is not covered by $Q_1 \cup \cdots \cup P_i \cup Q_{i+1}$. Let $m_i$ be upper left corner of $P_{i+1}$. That defines two half lines starting at $m_i$, one going down and one going to the right ($f_{i+1}$). They hit $S$ at two positions. They give us two other corner of $P_{i+1}$. This completes the definition of $P_{i+1}$.

The procedure stops when the already constructed rectangles cover all the $V_i$’s (or all the $h_i$’s).

Now we prove that the constructed system of rectangles satisfy (1)–(3).

(3) is immediate.

In order to prove (1) we need a few remarks.

It is immediate from the definition that as $i$ is increasing the lines, $e_i$’s are moving to the left and the lines $f_i$ are moving up.

The definition also implies that the upper left corner of $P_i$ is on $e_i$ or is left from $e_i$. Similarly the lower right corner of $Q_{i+1}$ is on $f_i$ or is below $f_i$. This guarantees that $Q_1 \cup \cdots \cup P_i \cup Q_{i+1}$ covers everything left from $e_{i+1}$ in the region between $s$ and $S$. Similarly $Q_1 \cup \cdots \cup P_i \cup Q_{i+1} \cup P_{i+1}$ covers everything below $f_{i+1}$ in the region $s$ and $S$. This proves (1).

(2) From the definition the top side of $Q_i$ (and this way the whole rectangle) is not above $f_i$. The lower right corner of $Q_{i+1}$ (let us say $M_{ij}$) is not above $f_i$, but it is the last maximal 1 with this property. This guarantees that $Q_{i+1}$’s lower right corner (and this way the whole rectangle) if above $f_i$. So the rows of $Q_i$ and $Q_{i+2}$ are completely disjoint. One gets the corresponding statements for the columns and for the $P_i$’s similarly. (2) is an easy consequence of this.

This completes the proof. \square

**Corollary 6.2.** (1) $f(n, m; C_{10})$ is linear.

(2) $f(n, m; \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix})$ is linear.

(3) $f(n, m; \left\{ \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix} \right\})$ is linear.
Proof. (1) Take the cover guaranteed by Theorem 6.1. Count the 1’s separately in different covering rectangles. We know that in the upper left corner or in the lower right corner there is a 1. So we can bound the number of 1’s using that

\[
f(n, m; \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}) \quad \text{and} \quad f(n, m; \begin{pmatrix} 1 \\ 1 \end{pmatrix})
\]

are linear. If we add up these bounds we obtain the claim in (1).

The same proof works for (2), but there we use Theorem 5.1. (3) follows the same way. □

7. Davenport–Schinzel matrices

In this and the next section we consider the complexity of

\[C_6 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}.
\]

Definition 7.1. A matrix \(M\) is called a Davenport–Schinzel matrix if it does not have \(C_6\) as a subconfiguration.

The naming is based on the analogy between this kind of matrices and Davenport–Schinzel sequences (see [7]).

The main result in this section is to construct a Davenport–Schinzel matrix with \(\Omega(n\alpha(n))\) many 1's. Finally we discuss other configurations, missing from our matrix.

Our construction is very similar to known constructions of Davenport–Schinzel sequences (see [12, 18]). We use the same double induction. But instead of sequences we work with matrices.

The matrices we are constructing have two parameters \(s\) and \(t\). We refer to them as \(M(s, t)\). First we describe a few properties of \(M(s, t)\). The recursive definition of these matrices is assuming these properties so we need to maintain them.

(a) The size of the matrix is \(tC(s, t) \times tC(s, t)\), where \(C(s, t)\) is defined as follows. \(C(s, t) = C(s, t-1)C(s-1, C(s, t-1))\) and \(C(1, s) = 1\) and \(C(s, 1) = 2\), for \(s > 1\).

(b) The \(tC(s, t)\) many rows are divided into blocks. We will refer to them as horizontal blocks. One block contains \(t\) rows (hence we have \(C(s, t)\) many blocks). Let \(H_i\) be the set of the \(((i-1)t+1)st, \ldots, (it)th\) rows, i.e., the \(i\)th horizontal block.

(c) Inside \(H_i\) the appearance of the first 1 happens in the same column (considering different rows). Let us say this is the \((c_i)th\) column. The 1's in these columns are called leading 1’s.
(d) 1 = c_1 < c_2 < c_3 < \cdots < c_{C(s,t)}$. These columns divide the matrix into vertical blocks. Let $V_i$ be the set of columns from the $(c_i)$th, through $(c_{i+1}-1)$st, i.e., the $i$th vertical block.

The definition of $M(s,t)$ is going to use the matrices $S = M(s,t-1)$ and $B = M(s-1,C(s,t-1))$. (Think about $S$ as a small matrix and about $B$ as a big matrix.) $B$ has $C(s-1,C(s,t-1))$ many horizontal blocks of size $C(s,t-1)$. $B$ has $C(s-1,C(s,t-1))$ many vertical blocks too. Let $v_i$ be the number of columns contained in the $i$th one. $S$ has $C(s,t-1)$ many blocks (one for each row in a block of $B$).

The following definition assumes properties (a)-(d). (So one must check that these properties are maintained.)

**Definition 7.2.** $M(1,s)$ is an identity matrix of size $s \times s$. $M(s,1)$ is

\[
\begin{pmatrix}
1 & 1 \\
1 & 0
\end{pmatrix}
\]

(for $s > 1$). In order to define $M(s,t)$ take $C(s-1,C(s,t-1))$ many copies of $S$ (one for each horizontal block of $B$). The construction of $M$ will be completed in $C(s-1,C(s,t-1))$ many stages. In the $i$th stage we add $(t-1)C(s,t-1) + C(s,t-1)$ many new rows and $(t-1)C(s,t-1) + v_i$ many new columns to the part already built. The construction starts with the empty matrix. The general (ith) stage is the following:

1. We put $(t-1)C(s,t-1)$ many new rows and new columns after the already existing ones. In the intersection of the new rows and columns we place a copy of $S$.
2. We insert an extra row after each horizontal block of the new copy of $S$. In these extra rows we place one extra 1, under each leading column.
3. Finally we add $v_i$ new columns (after the old ones). In the new space we place a copy of the $i$th vertical block of $B$ using the extra rows.

The constructed matrix $M = M(s,t)$ has properties (a)-(d). Let us introduce a few notations. *Ordinary rows* and *ordinary columns* are the rows and columns introduced in step (1). *Extra rows* are the rows introduced in step (2). *Extra columns* are the ones introduced in step (3). The 1's introduced in step (1) are the *ordinary* 1's. The 1's introduced in step (2) are called the *extra* 1's. The 1's introduced in step (3) are the *new* 1's.

The previous notations give a partition of 1's into new, ordinary and extra 1's. There are similar partitions for rows and columns.

Any extra 1 is in an ordinary column and in an extra row.

The next lemma summarizes a few simple statements about the matrix $M(s,t)$.

**Lemma 7.3.** (1) If $s$ and $t$ are chosen appropriately and $n = sC(s,t)$ then $M(s,t)$ is an $n \times n$ matrix with $n\alpha(n)$ many 1's.
(2) The \((c)\)th column contains 1’s inside \(H_i\) and no other 1’s.

(3) Inside \(H_i\), after the leading column the 1’s are decreasing, i.e., if \(k\) and \(l\) are two 1’s in the same horizontal block and they are not leading 1’s then \(k \uparrow \downarrow l\) or \(l \downarrow \uparrow k\). (Recall that \(q \uparrow \downarrow p\) vaguely means that \(p\) is in south, east or south-east direction from \(q\).)

(4) If \(l\) is a new 1 and \(k\) is a 1 such that \(l \uparrow \downarrow k\) then \(k\) is a new 1 too. (Recall that \(q \uparrow \downarrow p\) vaguely means \(p\) is in north, east or north-east direction from \(q\).)

(5) If \(l\) is an ordinary 1 and \(k\) is a 1 in \(l\)’s column or in \(l\)’s row then \(k\) is an ordinary 1 in the same horizontal block with the one exception when \(l\) is a leading 1 and \(k\) is the extra 1 in its column.

(6) If \(l\) is an extra 1 or an ordinary 1 and \(k\) is an ordinary 1 such that \(l \uparrow \downarrow k\) then \(l\) and \(k\) are in the same horizontal block.

**Proof.** For (1) we refer the reader to [12, 18].

The proof of (2)–(6) is easy induction following the definition of \(M(s, t)\). \(\square\)

Now we are ready to discuss the missing configurations in \(M(s, t)\).

**Theorem 7.4.** \(M(s, t)\) does not have the following configurations:

(i) \[
\begin{pmatrix}
1 \\
1 \\
1 \\
1
\end{pmatrix},
\]
(ii) \[
\begin{pmatrix}
1 & 1 \\
1 \\
1
\end{pmatrix},
\]
(iii) \[
\begin{pmatrix}
1 \\
1 \\
1
\end{pmatrix},
\]
(iv) \[
\begin{pmatrix}
1 & 1 \\
1 & 1 \\
1
\end{pmatrix},
\]
(v) \[
\begin{pmatrix}
1 & 1 \\
1 \\
1
\end{pmatrix},
\]
(vi) \[
\begin{pmatrix}
1 & 1 \\
1 \\
1
\end{pmatrix},
\]
(vii) \[
\begin{pmatrix}
1 & 1 \\
1 \\
1
\end{pmatrix},
\]
(viii) \[
\begin{pmatrix}
1 & 1 \\
1
\end{pmatrix}.
\]

**Proof.** Each configuration in the statement has four 1’s in it. Let us order these 1’s. A 1 is earlier than another if its row is earlier or if they are in the same row and it is left from the other. In the case of each configuration name the four 1’s as \(a, b, c\) and \(d\) following the previously defined order.

Our proof is by induction following the definition of \(M(s, t)\). The initial case is \(s = 1\) or \(t = 1\). Then the statement is clear.

The induction step is proved by contradiction. Let us assume that in \(M(s, t)\) we can find four different 1’s: the image of \(a, b, c\) and \(d\), such that they obey the configuration. The individual configurations are considered separately.
(i) We distinguish cases depending on what kind of entry corresponds to $c$. From now on we do not distinguish $a, b, c, d$ and their images.

Case 1: $c$ is an extra 1.

Then $d$ is a new 1. $a$ is in a leading column but it is not an extra 1. So $a$'s row is an ordinary row. On the other hand $c \nmid^* b$, hence (by Lemma 7.3 (4)) $b$ is a new 1. So $b$'s row (which is the same as $a$'s row) is an extra row. Contradiction.

Case 2: $c$ is a new 1.

Using Lemma 7.3 (4) the whole configuration consists of new 1's. So it can be recognized inside $M(s - 1, C(s, t - 1))$. Contradiction with the inductive hypothesis.

Case 3: $c$ is an ordinary 1.

Using Lemma 7.3 (5) the whole configuration consists of ordinary 1's from the same horizontal block. So our configuration can be recognized in a copy of $M(s, t - 1)$.

(ii) Case 1: $c$ is an extra 1.

Then $a, b$ and $d$ are new 1's. Let $c'$ the first 1 after $c$ in its row (that row is an extra row and $c'$ is a new 1). Easy to check that $a, b, c'$ and $d$ give us a configuration $C_1$ or one which is the same as the original configuration. So using (i) or the inductive hypothesis we get a contradiction.

Case 2: $c$ is a new 1.

$a, b, c$ and $d$ are all new 1's. So our configuration is in a copy of $M(s - 1, C(s, t - 1))$.

Case 3: $c$ is an ordinary 1.

Using Lemma 7.3 (5) our configuration is inside a copy of $M(s, t - 1)$.

(iii)-(vi) Using the same case analysis based on the bottom left 1 (which is not necessarily $c$).

(vii) Case 1: $d$ is an extra 1.

Then $a, b$ and $c$ are ordinary 1's in the same horizontal block (using Lemma 7.3 (5) and the fact that ordinary columns and rows in the same block are consecutive ones). Then the positions of $b$ and $c$ are contradictory with Lemma 7.3 (3).

Case 2: $d$ is a new 1.

The same as the previous second cases.

Case 3: $d$ is an ordinary 1.

The same as the previous third cases.

(viii) Case 1: $d$ is an extra 1 and $c$ is a new 1.

$c \nmid^* b$ hence $b$ is a new 1 too, in particular $a$'s and $b$'s row is an extra row. $d \nmid^* a$ so $a$ cannot be an extra 1. Hence all four 1's are new except $d$. Move $d$ right to the first 1 in its row. Then we obtain four new ones (hence they are in a copy of $M(s - 1, C(s, t - 1))$ such that their configuration is the one described in (vii) or in (viii).

Case 2: $d$ is an extra 1 and $c$ is an ordinary 1.

Using similar arguments as before we have that all four 1's are ordinary except $d$ and they are in the same horizontal block. Move $d$ up by one position. We
obtain four ordinary 1's (inside a copy of $M(s, t-1)$) such that their configuration is the one described in (vi) or in (viii).

Case 3: $d$ is not an extra 1.

In this case take the bottom left 1 ($d$) and replace it with another 1 by shifting it to the lending 1 in its row and sinking it to the bottom 1 in that column. This way we obtain the same configuration but the new $d$ is an extra 1. This was handled in the previous cases. □

The previous theorem gives lower bounds on the complexity of several configurations and sets of configurations.

**Corollary 7.5.** (1) $f(n; C_5) = \Omega(n\alpha(n))$, (2) $f(n; C_8) = \Omega(n\alpha(n))$,

(3) $f(n; \{ \begin{pmatrix} 1 & 1 \\ 1 & 1 \\ 1 & 1 \\ 1 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 1 \\ 1 & 1 \\ 1 & 1 \\ 1 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 1 \\ 1 & 1 \\ 1 & 1 \\ 1 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 1 \\ 1 & 1 \\ 1 & 1 \\ 1 & 1 \end{pmatrix} \} = \Omega(n\alpha(n))$.

8. Upper bound on Davenport–Schinzel matrices

In this section we prove that

**Theorem 8.1.** $f(n; \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}) \leq \mathcal{O}(n\alpha(n))$.

**Proof.** Let $A' = (a_{ij})$ be an $n \times m$ 0–1 matrix not having a subconfiguration of $C_6 = \begin{pmatrix} 1 & 1 \\ 1 & 1 \\ 1 & 1 \end{pmatrix}$. Delete the first and the last 1 in each row, and keep only the columns with at least 2 entries. The obtained matrix is denoted by $A = (a_{ij})$, and obviously, for the number of entries we have $\|A'\| \leq \|A\| + 2n + m$. Form a sequence $s_j$ from the $j$th column $(a_{ij})_{1 \leq i \leq \ell} =: l(j)$ in the following way

$$s_j = (s_1, s_2, \ldots, s_{l(j)}),$$

where $s_1 < s_2 < \cdots < s_{l(j)}$ and $a_{ij} = 1$ for $1 \leq i \leq l(j)$. Form one sequence $s' =: s_1s_2\cdots s_m$ in this order. Delete from $s'$ the element $s_{l(j)}$ if it equals to $s_{l(j)}^{-1}$. In the obtained sequence $s$ there are no equal consecutive elements. We claim that $s$ does not contain a subsequence $ababa$, i.e., it is a $DS_6(n)$ sequence.

Suppose on the contrary. Then there exists a subsequence $abab$ of $s$ with $a < b$. So there are $j_1 \leq \cdots \leq j_4$ such that $a \in s_{j_1}$, $b \in s_{j_2}$, $a \in s_{j_3}$, $b \in s_{j_4}$. Here $j_2 < j_3$. Otherwise the first $b$ in $abab$ could not precede the second $a$ in $s$. Consider the submatrix defined by the rows $a$ and $b$ and the columns $\{j_1, \ldots, j_4\}$. There are
four possibilities.

\[ j_1 < j_2 < j_3 < j_4, \quad \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}, \]

\[ j_1 = j_2 < j_3 < j_4, \quad \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}, \]

\[ j_1 = j_2 \text{ and } j_3 = j_4, \quad \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}, \]

\[ j_1 < j_2 < j_3 = j_4, \quad \begin{pmatrix} 1 & 1 \\ 1 & 1 \end{pmatrix}. \]

In each case \( A \) will contain a \( C_\alpha \), a contradiction.
So \( \|A\| \leq 2n + 2m + ds_3(n) \).

**Corollary 8.2.** \( f(n; C_7), f(n; C_9) = O(n \alpha(n)) \).

In the very same way we can obtain the following theorem. Let \( C^{2k} \) be a partial \( 2 \times 2k \) matrix with \( c_{1,2i-1} = 1, \ c_{2,2i} = 1 \) for \( 1 \leq i \leq k \).

**Theorem 8.3.** \( f(n; C^{2k}) = O(ds_{2k-3}(n)) \).

It is not difficult to give a lower bound for \( f(n; C^{2k}) \) which is probably closer to \( f \) as the upper bound.

**Theorem 8.4.** \( f(n; C^{2k}) > \Omega(ds_{2k+1}(n)/ \alpha(\alpha(n))^O(\alpha(\alpha(n))^{2k-4})) \).

**Remark.** Here the right-hand side is superlinear. For the best bound on \( ds_{2k+1}(n) \) see [3].

**Proof of Theorem 8.4.** Let \( s \) be a Davenport–Schinzel sequence \( s \in DS_{2k+1}(n) \) of length \( ds_{2k+1}(n) \) such that the element \( i \) appears earlier than \( j \) for \( i < j \). It is well known [16], that

\[ ds_{2k+1}(n) = O(n^2 \alpha(n)^{O(\alpha(n))^{2k-4}}) \]

Split \( s \) into \( n \) almost equal parts \( s = s_1 \cdots s_n \). Let \( s'_i \) be a set of distinct values of \( s_i, \|s'_i\| = x \). We have that

\[ y = \frac{ds_{2k+1}(n)}{n} - 1 \leq \|s'_i\| \leq xO(\alpha(x)^{O(\alpha(x))^{2k-4}}) \leq xO(\alpha(y)^{O(\alpha(y))^{2k-4}}). \]

Here \( y = O(\alpha(n)^{O(\alpha(n))^{2k-4}}) \), so \( \alpha(y) = \alpha(\alpha(n)) + O(1) \).

Finally, forming the \( i \)-th column of an \( n \times n \) matrix \( A \) from \( s'_i \) we obtain the desired configuration without \( C^{2k} \).
9. Conclusions and open problems

Table 1 summarizes our results.

Finally we mention several open problems. The first few ones are suggested by the table. Even in the case of configurations with four 1's there are several unknown complexities.

Is it true that the complexity of all permutation configurations are linear?

What is the characterization of configurations with linear complexity? In extremal graph theory the forbidden subgraphs with linear threshold are exactly the trees.

Is it true, that if $G$ is the (bipartite) graph corresponding the configuration $C$ then

$$f(n; C) < O(T(n; G) \log n)?$$

(9.1)

Does (9.1) hold at least for trees?

There are several combinatorial structures with an underlying order where the

<table>
<thead>
<tr>
<th>Configurations</th>
<th>Lower bound</th>
<th>Upper bound</th>
</tr>
</thead>
<tbody>
<tr>
<td>((1 1 1 1)^\frac{1}{2})</td>
<td>(\Theta(n^{3/2}))</td>
<td>(\Theta(n^{3/2}))</td>
</tr>
<tr>
<td></td>
<td>see Theorem 1.1</td>
<td></td>
</tr>
<tr>
<td>((1 1 1 1)^\frac{1}{2}), ((1 1 1 1)^\frac{1}{2})</td>
<td>(\Theta(n \log n))</td>
<td>(\Theta(n \log n))</td>
</tr>
<tr>
<td></td>
<td>see Theorem 3.1 and Corollary 3.4</td>
<td></td>
</tr>
<tr>
<td>((1 1 1 1)^\frac{1}{2})</td>
<td>(\Omega(n \log n))</td>
<td>(O(n \log n))</td>
</tr>
<tr>
<td></td>
<td>see Corollary 4.3</td>
<td></td>
</tr>
<tr>
<td>((1 1 1 1)^\frac{1}{2})</td>
<td>(\Omega(n))</td>
<td>(O(n \log n))</td>
</tr>
<tr>
<td></td>
<td>see Corollary 3.4</td>
<td></td>
</tr>
<tr>
<td>((1 1 1 1)^\frac{1}{2}), ((1 1 1 1)^\frac{1}{2})</td>
<td>(\Theta(n \alpha(n)))</td>
<td>(\Theta(n \alpha(n)))</td>
</tr>
<tr>
<td></td>
<td>see Corollary 7.5 and Theorem 8.1</td>
<td></td>
</tr>
<tr>
<td>((1 1 1 1)^\frac{1}{2}), ((1 1 1 1)^\frac{1}{2})</td>
<td>(\Omega(n))</td>
<td>(O(n \alpha(n)))</td>
</tr>
<tr>
<td></td>
<td>see Corollary 8.2</td>
<td></td>
</tr>
<tr>
<td>All the other 28 matrices with 4 entries</td>
<td>(\Theta(n))</td>
<td>(\Theta(n))</td>
</tr>
</tbody>
</table>
similar extremal question is interesting. An example is a set of intervals on a given line. How many intervals (over \( n \) endpoints) guarantee the existence of a given interval configuration? Similar questions can be asked about diagonals in a cycle. Davenport and Schinzel's original question can be extended to arbitrary forbidden subsequences. As far we know there is no organized account of these questions.

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References