A subexponential lower bound for branching programs 
restricted with regard to some semantic aspects

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Abstract

Branching programs (b.p.s) or binary decision diagrams are a general graph-based model of sequential computation. The b.p.s of polynomial size are a nonuniform counterpart of LOG. Lower bounds of different kinds of restricted b.p.s are intensively investigated. The restrictions based on the number of tests of variables during any computation (along any path in the case of syntactic b.p.s, resp.) are very important. Various superpolynomial or even exponential lower bounds are known for 1-branching programs, (1, +k)-branching programs, and syntactic k-branching programs.

In the paper we present a restriction of another type - the so called gentle branching programs (g-b.p.s) together with the following results.

1. For each 1-b.p. \( P \) there is a gentle 1-b.p. \( P' \) computing the same function such that \( |P'| \leq 3|P| \).
2. Some Boolean functions which are superpolynomially (or even exponentially) hard for 1-b.p.s are polynomially easy for g-b.p.s.
3. The same holds for Boolean functions which are superpolynomially hard for (1, +k)-b.p.s or for syntactic k-b.p.s.
4. The functions whose sequences cause the hierarchies for (1, +k)-b.p.s and for syntactic k-b.p.s (both with respect to \( k \)) are polynomially easy for g-b.p.s.
5. We find a function in \( \mathcal{P} \) which is superpolynomially hard for g-b.p.s.

The proof is based on a lower bound theorem.

Both, the lower bound theorem and the definition of gentle branching programs are derived from a deeper consideration of the phenomenon of branching of computations.

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1 Introduction

1.1 Branching programs

A branching program (b.p.) is a computational model for representing Boolean functions. The input of a branching program is a vector consisting of the values of \( n \) Boolean variables. The branching program itself is a directed acyclic graph with one source. The out-degree of each node is at most 2. Every branching node, i.e., a node of out-degree 2, is labeled by an input variable (or equivalently: by an input bit) and one of its out-going edges is labeled by 0, the other one by 1. The sinks (out-degree = 0) are labeled by 0 and 1. A branching program determines a Boolean function as follows. The computation starts at the source. If a node of out-degree 1 is reached, the computation follows the unique edge leaving the node. In each branching node, the variable assigned to the node is tested and the out-going edge labeled by the actual value of the variable is chosen. Finally, a sink is reached. Its label determines the value of the function for the given input. By the size of a branching program \( P \) we mean the number \(|P|\) of its nodes.

Branching programs are a model of the configuration space of Turing machines where each node corresponds to a configuration. Thus, the polynomially sized b.p.s represent a nonuniform variant of \( \text{LOG} \). Hence a superpolynomial lower bound on b.p.s for a Boolean function computable within polynomial time would imply \( P \neq \text{LOG} \).

Some generalization of b.p.s have been studied. Let us mention nondeterministic b.p.s where each branching node may have more than two out-going edges. The polynomial sized nondeterministic b.p.s represent a nonuniform variant of \( \text{NLOG} \). Let us mention also the \( \omega \)-b.p.s [13], and the probabilistic b.p.s [1].

In the present paper we use the classical deterministic branching programs. We work in one of the main streams of branching programs investigation which consists in proving lower bounds for restricted branching programs, it means for branching programs of some special properties.

An example of a restriction from the early beginning of investigation are branching programs of restricted width [3] or oblivious branching programs of limited depth. [11].

A special type of b.p.s are the so called ordered binary decision diagrams (OBDD) where along each computation the tests preserve a fixed ordering of variables and each variable is tested at most once. OBDD’s play an important role as a tool for representing Boolean functions in many computer applications.
1.2 Restricted branching programs

The most commonly used restrictions are based on the number of tests during a computation or along a path in b.p.s in the question. So, the read-once branching programs (1-b.p.s), where during each computation each variable is tested at most once, are well-known. The lower bounds range from $2^{\sqrt{n}}$[23][24][5] through $2^{n/\log n}$ [8] to $2^n$ [12][22] and $2^{n-o(n)}$ [19]. There are lower bounds also in [17][1],[7].

There are some attempts to prove lower bounds for large classes of restricted branching programs. The real-time branching programs are allowed to do at most $n$ tests during each computation. The related lower bounds are in [6],[25],[12].

The syntactic $k$-branching programs are allowed to have at most $k$ tests on each variable along any (consistent or nonconsistent) path. The lower bounds are in [14],[4],[9] and the hierarchy with respect to $k$ in [15],[16].

The $(1,+k)$-branching programs are allowed to test at most $k$ variables repeatedly during each computation. The lower bounds are in [26][18][10] and the hierarchy in [20]. (The hierarchy for the syntactic case is proven in [21].)

We present another attempt in this direction with results as stated in the Abstract. However, the problem of the superpolynomial lower bound for 2-b.p.s (from 1983) remains open till now.

1.3 The basic ideas and the results

Our results are derived from two informal ideas.

Let us introduce the first one. For the branching program in question given any input we want to catch what information about the contents of the input bits is remembered (or forgotten) at any moment of the computation. We represent such a knowledge by a word of the length $n$ over the alphabet \{0, 1, +, #\} where +, # stand for "unknown".

From some reasoning about TM’s it follows that the definition of crosses (+) has to satisfy the next two, very natural conditions:

a) Immediately before a test on a bit $i$, $i$ is crossed.

b) Immediately after a test on a bit $i$, $i$ is non-crossed.

From this requirement it follows that the knowledge depends not only on the node reached during the computation but that different inputs reaching the same node may have remembered different information. This notion of remembered and forgotten information formally defined in Section 3 enables us to formulate and to prove our two lower bound method theorems.

The second informal idea concerns a situation where two computations reach
the same node \( v \) of the b.p. and after \( v \) they never branch. In this case we see that the bits on which these inputs differ have no influence on the common part of their computations and hence no influence on the reached sink. At \( v \) these computations must have achieved a partial result. In Section 3 we try to catch this strange phenomenon by defining when a bit of an input has \(#\). This allows us to define the distribution of inputs over any branching program which is a key-notion of the definition of gentle branching programs in Section 4. There we also prove the relation of 1-b.p.s and gentle b.p.s as indicated in the Abstract.

In Section 5 we prove a subexponential lower bound \( 2^{n^{1/8}} \) on the size of gentle branching programs for a Boolean function which is computable within polynomial time on Turing machines. The proof is based on the theorems from Section 3 mentioned above.

In Section 6 we prove an \( O(n) \) upper bound (on gentle branching programs) for the Boolean function \( f \) which is defined as follows: \( f(A\bar{x}) = 1 \) iff \( A\bar{x} = \emptyset \) (the input bits are arranged into a matrix and a vector). From [9][18] it is possible to derive that \( f \) is subexponentially hard for syntactic \( k \)- and for \((1,+k)\)-branching programs. Moreover we prove an \( O(n\log n) \) upper bound for functions which are used in the Okolnishnikov extension of the hierarchy of syntactic \( k \)-branching programs [15]. Further, we prove a polynomial upper bound \( O(n^2) \) for multipointer functions which are used in [20] for the proof of the hierarchy of \((1,+k)\)-branching programs.

In Section 7 we prove an \( O(n^{3/2}) \) upper bound for the parity of the number of triangles in a graph (this function requires \( 2^{\Theta(n)} \) on read-once b.p.s [2][22]).

Further, we prove a \( O(n^2) \) upper bound for half-cliques-only function and a \( O(n^3) \) upper bound for the Ablayev-Karpinski function [1]. The last upper bound is proven for the multiplication. Its size is \( O(n^3) \) (cf. [17]).

## 2 Preliminaries

First we introduce some technicalities concerning words. For a binary word (a binary string) \( m \in \{0,1\}^n \) by \( m_i \) we mean its \( i \)-th symbol. If \( A \subseteq \{1,2,...,n\} \), by \( m \upharpoonright A \) we mean the assignment of bits from \( A \) consistent with \( m \). If \( R \) is a predicate on \( \{1,2,...,n\} \) by \( m_R \) we mean \( m \upharpoonright A \) where \( A = \{ j \mid R(j) \} \). E.g. \( m_{\neq i}, m_{<i} \) where \( 1 \leq i \leq n \). If \( A, B \subseteq \{1,...,n\}, A \cap B = \emptyset \) and \( a, b, \text{resp.} \) is an assignment of \( A \) (\( B, \text{resp.} \) then by \([ab] \) we mean the assignment of \( A \cup B \) such that \([ab] \upharpoonright A = a, [ab] \upharpoonright B = b \). E.g. \([m_{\leq i}\,m_{>i}] = m \). For a word \( m \) by \(|m|_1 \) we mean the number of its ones; similarly \(|m|_0 \).

Secondly we introduce some notions concerning branching programs. Let \( P \) be a branching program, \( m \in \{0,1\}^n \) be an input word. By \( \text{comp}(m) \) (the computation on \( m \)) we understand the sequence of nodes of \( P \) which starts in the source of \( P \) and ends in a sink such that in each branching node \( v \) with a
label \( i \) \( \text{comp}(m) \) chooses the successor \( v' \) such that the label of the edge \((v, v')\)

is equal to \( m_i \). If a node \( v \in \text{comp}(m) \) we say that \( \text{comp}(m) \) goes through \( v \)
or simply that \( m \) goes through \( v \) or that \( m \) reaches \( v \). Similarly we say that \( m \)
go through an edge.

If \( v \) is a node of a branching program \( P \) then by \( P_v \) we mean the branching

program with the source in \( v \). If \( M \subseteq \{0, 1\}^n \) then by \( P \downharpoonright M \) we mean the

branching program which is given by deleting of all nodes and edges (from \( P \))

for which no \( m \in M \) goes through.

Further, let us remember the operation of development of a branching pro-

gram \( P \) (from its source) into a tree. By a development of \( P \) from its node \( v \)

we understand the development of \( P_v \).

By a subcube we mean the set \( M \) of inputs such that there is \( I, I \subseteq \{1, ..., n \} \)

and \( M = \{ |u| | u \in \{0, 1\}^I \} \) where \( u \) is an assignment of \( \{1, ..., n \} \) \( - I \). The bits

from \( I \) are called free bits, the other are called fixed bits. For a sink \( v \) of \( P \)

\( M \) is called an original subcube if all inputs from \( M \) reach \( v \) and the free bits of

\( M \) are never tested. For a node \( v \in P \) by an aggregate computation we mean

the set of computations (the set of inputs) which follow the same branch of \( P_v \)
till its sink.

3 \hspace{1em} \textbf{Windows}

By a \textit{window} we mean a string \( w = w_1w_2...w_n \in \{0, 1, +, \#\}^n \). By the length \( lw \)
of such a window \( w \) we understand the number of i’s such that \( w_i \in \{0, 1\} \) (the

number of noncrossed bits).

We say that a window \( w \) is a window over a word \( u \in \{0, 1\}^n \) iff for each its

noncrossed bit \( i \) \( w_i = u_i \) holds.

\begin{definition}
Let \( P \) be a branching program (over \( n \) variables). Let \( v \) be a

node of \( P \) and \( M \) be a subset of the set of all inputs which (starting in the

source of \( P \)) reach \( v \). Let \( m \in M \).

Then we define the window \( \omega(m, v, M, P) \) over \( m \) in the node \( v \) with respect

to \( M \) as follows: (Starting at \( v \) we develop the program \( P_v \) into a tree. In this

tree we perform all computations starting from \( v \) which are given by the inputs

from \( M \). We omit all the edges and the nodes of the tree which are not reached

by any of these computations. We omit all non-branching nodes. Now let us

consider the branch \( b \) which is followed by the computation \( \text{comp}(m) \).

1. We assign crosses (+) to all bits of \( \omega(m, v, M, P) \) which are labels of

nodes of the branch \( b \).

2. Now we consider the set \( L \) of all inputs (from \( M \)) which follow \( b \) till

its leaf. We assign the cross (+) to each bit which is not tested during any

computation (from the source to the sink) on any word from \( L \).

3. Further, we assign a double-cross \( \# \) to each remaining bit \( i \) of

\( \omega(m, v, M, P) \) for which in \( L \) there is an input \( m' \in M \) such that \( m_i \neq m'_i \).

4. The remaining bits of \( \omega(m, v, M, P) \) have the same contents as in \( m \).

\end{definition}
**Comment.** a) We see that all inputs following the same branch $b$ have the same set of crossed (+) bits, of double-crossed (#) bits and that all non-crossed bits have the same contents.

b) The larger $M$ the lesser number of non-crossed bits.

Before stating the main theorems we will prove some facts to become familiar with the definition.

For a moment due to some technical reasons we shall use a slightly modified branching programs such that on each path for each node with in-degree $>1$ or out-degree $>1$ both its immediate predecessors and successors have in-degree = out-degree = 1.

Let us fix a set $M \subseteq \{0,1\}^n$, let $P$ be a branching program such that $P = P \upharpoonright M$. For each node $v$ of $P$ we put $M_v = \{ m \in M | m \text{ reaches } v \}$. In the sequel (before the main theorems) at each node $v \in P$ we will consider the windows with respect to $M_v$ only.

**Proposition 3.2** Let $v$ be a branching node of $P$ with a test on the bit $i$. Let $v_1$ be one of two nodes which immediately succeeds $v$. Then for any $m \in M$ going through $v$, $v_1$ the following holds: The windows for $m$ at $v$ and at $v_1$ differ only in the bit $i$. At $v$ $i$ is crossed, at $v_1$ $i$ is non-crossed.

**Proof:** Let $m \in M$. Let $b$ be the branch of the tree developed at $v$ which is followed by $m$. Since $P = P \upharpoonright M$ at $v$ the tree has a branching node labeled by $i$. Hence $i$ is crossed. At $v_1$ all inputs from $M_{v_1}$ have the same value of $i$, hence $i$ is non-crossed. The branch $b_1$ of the tree developed at $v_1$ (induced by $m$) has the same branching nodes as $b$ has (with the exception of $v$ with the label $i$) and its leaf has the same set of inputs as the leaf of $b$ has. Hence the set of crossed and double-crossed bits remains the same (with the exception of $i$). □

**Comment.** We can follow a computation on any input $m$ from the source to a sink and consider the changes of the window on $m$ during the computation. The proposition says that a test on a bit implies that exactly one crossed bit becomes non-crossed, which is in a good correspondence with the informal idea that a test is an acquisition of exactly one bit of information.

By a symmetric word we mean any word of the form $wv$ where $v = u^R$. By a pair of symmetric positions we mean any pair $(i, j)$, $1 \leq i, j \leq n$, where $i = n - j + 1$.

**Theorem 3.3** Let $M = \{0,1\}^n$, let $P$ be a branching program recognizing symmetric words of the length $n$. Then for each symmetric word $m$ and for each pair $(i, j)$ of symmetric positions during $\text{comp}(m)$ there is at least one node $v$ at which, in the window for $m$ (with the respect to $M_v$), both $i, j$ are non-crossed.
Comment. The theorem corresponds to the informal idea that for comparing two bits $i, j$ it is necessary that both $i, j$ are known (non-crossed) at the same moment.

Proof:

By contradiction. Let $m$ be a symmetric word and $i, j$ be a pair of symmetric positions such that during $\text{comp}(m)$ they are never both non-crossed at the same moment (at the same node). We see that for both $i, j$ there are some nodes at which they are non-crossed since both must be tested during $\text{comp}(m)$. Let us assume that the last window on $i$ precedes the last window on $j$. Let $v$ be the node from $\text{comp}(m)$ such that immediately before $v$ is non-crossed and this is the last moment when $i$ is non-crossed during $\text{comp}(m)$. Since at $v$ is not non-crossed there is an input $m'$ such that $m_i \neq m_i'$ and such that a) $m, m'$ branches after $v$ on $i$ or b) $m, m'$ follow a common path till a common sink. According to a) there is a window on $m$ with the non-crossed $i$ after $v$. It is impossible. It remains only the case b). From b) it follows that $m_j'$ is also a symmetric word. After $v$ becomes non-crossed. Hence after $v$ there is a test on $j$ (according to the previous proposition). Since from $v$ $m, m'$ follow the same path $m, m'$ have the same value on $j$. For our symmetric words $m, m'$ we have $m_j' \neq m_i = m_j = m_j'$. A contradiction. □

Now we introduce two main lower bound method theorems.

**Theorem 3.4** Let $P$ be a branching program, let $v$ be one of its nodes, let $M$ be a subset of the set of all inputs which (starting at the source) reach $v$. Then the average number of crosses and double-crosses in $w(m, v, M, P)$ for $m \in M$ is at least $\log_2 |M|.

More formally: $\sum_{m \in M}(n - lw(m, v, M, P))/|M| \geq \log_2 |M|.$

Proof: Let us develop the program $P_v \downarrow M$ into a tree $T_1$. Let $d_1$ be the average number of crosses (+) and double-crosses (#) over all $m \in M$.

We transform $T_1$ into a tree $T_2$ which in each of its leaves has only one input from $M$. In each leaf of $T_1$ which is reached by two or more inputs from $M$ we start a new subtree such that on each its branch there is a test on each bit which is crossed according to rule 2 of the definition of $w(., v, M, P)$ and then on each branch there are some tests on bits which are double-crossed in $w(., v, M, P)$ until in each leaf of the transformed tree $T_2$ there is only one input from $M$.

If we apply the rules 1 - 3 from the definition of windows on $T_2$, we see that for each $m \in M$, for each bit $i$, if $i$ is crossed according to $T_2$ then $i$ was a crossed or double-crossed bit according to $T_1$. Since $T_2$ does not give double-crosses we see that for the average number $d_2$ of crosses given by $T_2$ the inequality $d_1 \geq d_2$ holds.

Further, we shall transform $T_2$ into a balanced tree $T_3$. (By a balanced tree we mean a tree where for each pair of its branches the difference of their lengths
is at most one.) We shall also transform the set $M$ into another set $M'$ of inputs such that $|M'| = |M|$. Let $d_3$ be the average number of crosses given by $T_3$ to inputs from $M'$. We shall prove $d_2 \geq d_3 \geq \log_2 |M'| = \log_2 |M|$.

Let us describe one step of the transformation of $T_2$ into $T_3$. Suppose there are two branches $b_1, b_2$ of $T_2$ such that $|b_1| - |b_2| \geq 2$ where for any branch $b \vdash b$ denotes the number of its tests (its length). Let $m_1$ ( $m_2$, resp.) be the input which follows $b_1$ ($b_2$, resp.) till its leaf. We delete $m_1$ from $M$. The total number of crosses is decreased by $|b_1| + 1$ ($|b_1| + 1$ is the number of crosses on $b_1$ and 1 is given by the disappearing of the last cross on the input $m$ which branches with $m_1$ in the last test of $b_1$). We add such an input $m'$ to $M$ that $m'$ follows $b_2$ till its leaf where $m'$ branches with $m_2$. By the last action the total number of crosses is increased by $|b_2| + 2$ ($|b_2| + 1$ is the number of crosses on $m'$ and 1 is for the additional cross on $m_2$).

We repeat this procedure until the resulting tree is balanced. We see that $|M| = |M'|$ and $d_2 \geq d_3$.

Now it remains to prove the following proposition.

**Proposition 3.5** Let $T_3$ be a ballanced tree. Let $M$ be the number of its leaves. Then the average length of branches of $T_3$ is at least $\log_2(|M|)$.

**Proof:** Let $k$ be the length of the shorter branches of $T_3$. Let $m$ be the number of nodes of $T_3$ on the level $k$ which are not leaves, $m \in <0, 2^k - 1 >$.

We see that the number of branches of $T_3$ is $2^k + m$. Hence, the average length of branches of $T_3$ is $[(2^k + m), k + 2m](/2^k + m)$. We have to verify $\log_2(2^k + m) \leq k + 2m/(2^k + m)$. It suffices to verify $2m/(2^k + m) \geq \log_2(1 + m/2^k)$. For $m \in <0, 2^k >$ we put $x = (1 + m/2^k)$; hence $x \in <1, 2 >$.

It suffices to verify $2x - 2 \geq x, \log_2 x$. It is the truth for $x = 1, x = 2$. Moreover for $x \in <1, 2 > (x, \log_2 x)^\prime\prime = (\log_2 e)/x > 0$. Hence, $x, \log_2 x$ is a convex function on $<1, 2 >$. Therefore $2x - 2 \geq x, \log_2 x$.

□

□

**Theorem 3.6** Let $P$ be a branching program. Let $\{X_i\}_{i=1}^r$ be a system of sets of some inputs. To each $X_i$ a node $v_i$ of $P$ is assigned such that all $x \in X_i$ go through $v_i$. Then

$$\log_2 r \geq \log_2 \sum_{i=1}^r |X_i| - n + \text{average length of } w(x, v_i, X_i, P).$$

More formally:

$$\log_2 r \geq \log_2 \sum_{i=1}^r |X_i| - n + \sum_{i=1}^r \sum_{x \in X_i} w(x, v_i, X_i, P)/ \sum_{i=1}^r |X_i|.$$
For \( x \in (0, \infty) \) we put \( \phi(x) = x, \log_2 x \). We know that \( \phi(x) \) is a convex function since \( \phi''(x) = (\log_2 e)/x > 0 \). According to Jensen's inequality we have
\[
\phi\left(\sum_{i=1}^r |X_i|/r\right) \leq \left(\sum_{i=1}^r \phi(|X_i|)\right)/r.
\]

Further \( \sum_{i=1}^r |X_i|, \log_2(\sum_{i=1}^r |X_i|/r) \leq \sum_{i=1}^r |X_i|, \log_2 |X_i| \) and
\[
\log_2 r \geq \log_2 \sum_{i=1}^r |X_i| - n + \text{average length of windows}.
\]

According to Theorem 3.4, the last expression is not greater than the average number of crosses and double-crosses in \( w(\cdot, v_i, X_i, P) \) which is equal to \( n - \text{average length of windows} \).

Hence
\[
\log_2 r \geq \log_2 \sum_{i=1}^r |X_i| - n + \text{average length of windows}.
\]

\( \square \)

Comments. a) Under the assumption that \( \{X_i\}_{i=1}^r \) cover all \( 2^n \) inputs the statement of the Theorem is abbreviated as follows:
"\( \log_2 r \geq \text{average length of windows} \)."

If there is a relation between \( r \) and \( |P| \) we have a lower bound for \( |P| \).

b) The Theorem corresponds to an informal idea that "If we want to have a possibility to remember much information about many inputs then we need a large memory".

c) Both Theorems form a method for proving lower bounds on the size of general branching programs. We will apply them in Section 5 for our lower bound. There is an open question to find another nontrivial application.

4 Gentle branching programs

Let us introduce the key-notation for the definition of the gentle branching programs.

Definition 4.1 By the distribution on a branching program we mean a partition of the set of all inputs \( \{0, 1\}^n \) together with a one-to-one assignment of all classes of the partition to some nodes of the branching program in question. The partition and the assignment are given by the following rule.

Briefly (with a danger of confusion): We stop each input in the node where its window has a double-cross \# for the first time (or in a sink).

Precisely: By Pr we mean the next procedure.

Input: \( A \subseteq \{0, 1\}^n \). For each computation \( \text{comp}(m) \), \( m \in A \), in each node \( v \) let us consider the window \( w(m, v, M, P) \) where \( M \subseteq A \) is the set of all inputs from \( A \) reaching \( v \).

For each \( m \in A \) we mark the first \( v \) in \( \text{comp}(m) \) such that \( w(m, v, M, P) \) has a bit with a double-cross \#. 

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In the set $S$ of marked nodes we say that $v \in S$ is a maximum one iff there is no $v_1 \in S$, $m_1 \in A$ such that $v_1$ is marked for $m_1$ and $v_1$ precedes $v$ in $\text{comp}(m_1)$.

For each maximum node $v \in S$ we define $M_v = \{ m \in A \mid v \text{ is marked for } \text{comp}(m) \}$.

Output: $R_A = \{ M_v \mid v \text{ is maximum} \}$.

$A := A - \bigcup R_A$.

We perform $Pr$ starting with $A = \{0,1\}^n$ until $R_A = \emptyset$.

The resulting partition is given as $R = \bigcup R_{A_1} \cup ... R_{A_k} \cup \{ M_{s_1}, \ldots M_{s_l} \}$ where $R_{A_i}$ means $R_A$ for the $i$-th cycle of $Pr$, and $M_{s_j}$ is the set of inputs which reach the sink $s_j$ (after the last cycle of $Pr$).

It is easy to see that the assignment $v \rightarrow M_v \in R$ which we have constructed is one-to-one.

By a proper class of the distribution we mean a class from $R_{A_1} \cup \ldots \cup R_{A_k}$.

It is clear that each improper class $M_{s_j}$ is an original subcube.

**Definition 4.2** For a distribution $\mathcal{D}$, by the factorization of $\mathcal{D}$ we mean a finer partition given by the factorization of each proper class of $\mathcal{D}$ according to the equivalence "to have the same set of (bits with) double-crosses #’s".

More precisely: Let $M$ be a proper class of $\mathcal{D}$ which is distributed at a node $v$. We say that $m_1, m_2 \in M$ are equivalent if $\forall i (i = 1 \ldots n)(w(m_1, v, M, P)_i = # \iff w(m_2, v, M, P)_i = #)$.

Let $M, C$ be two sets of inputs, $C \subseteq M$, and let all inputs from $M$ reach a node $v$. We say that $C$ preserves the double-crosses of $M$ if for each $u \in C$ and for each $i, 1 \leq i \leq n$ the next implication holds $w(u, v, M, P)_i = # \Rightarrow w(u, v, C, P)_i = #$.

It is easy to see that each class $F$ of the factorization of any class $M$ of the distribution preserves the double-crosses of $M$.

We say that a class of $M$ of a distribution $\mathcal{D}$ is a significant one iff $|M| \geq 2^n/\alpha$. We say that an $fd$-class $F$ subclass of a class $M$ of a distribution is a significant one iff $|F| \geq |M|/\beta$. $\alpha, \beta$ will be quantified in the next chapter.

**Definition 4.3** Let $F$ be an $fd$-class.

By $h_F$ we mean the set of bits which are double-crossed; by $t_F$ we mean the set of the other bits.

$H_F = \{ m \mid h_F \mid m \in F \}$.

$T_F = \{ m \mid t_F \mid m \in F \}$.

**Definition 4.4** Let $F$ be an $fd$-class at a node $v$. Let $D$ be the set of bits which are double-crossed (for all $f \in F$). Let $S_{F}$ be the tree induced by $F$ at $v$.
Let $V_1, \ldots, V_i$ be the aggregate computations which coincide with branches of $S_F$. Let $M_1^*, \ldots, M_i^*$ be the maximum original subcubes which are formed by inputs following $V_i$. By $G_i^j$ we mean the set of bits from $D$ which are free for $M_i^j$ and by $Q_i^j$ we mean the set of bits from $D$ which are fixed for $M_i^j$.

Definition 4.5 Let $P$ be a branching program. We say that $P$ is a gentle branching program if the following holds:

If there is a significant proper class in the distribution on $P$ then there is another significant proper class (at a node $v$) containing a significant fd-class $F$ fulfilling the following conditions (R1), (R2), (R3)

(R1) Let the sets $W_1, W_2$ of bits be defined as follows:

$W_1 = \{ \{ \forall f_1, f_2 \in F \mid w(f_1, v, F, P) = w(f_2, v, F, P), \in \{0, 1\}\}$

$W_2 = \{ \{ \exists f \in F \mid w(f, v, F, P), \in \{0, 1\}\}$. The condition is $|W_2 - W_1| \leq \gamma$. ($\gamma$ will be quantified in the next section.)

(R2) $\forall h \in H_F \forall t \in T_F \ f_P([h, t]) = f_Q(t)$ where $Q = P_v \mid F$.

(R3) For all $i, j \ G_i^j$ is the same set.

From the definition of windows we know that the condition (R3) implies for all $i, j \ G_i^j = \emptyset$.

Theorem 4.6 For each 1-branching program $P$ there is a gentle 1-branching program $P'$ with $|P'| \leq 3|P|$ and $f_P = f_{P'}$.

Proof: Let $P$ be a 1-branching program.

Proposition 4.7 Let $C$ be a proper class of the distribution on $P$. Then $C$ consists of a unique fd-subclass satisfying (R1), (R2) of the definition of gentle programs.

Proof:

Let $C$ be equal to the set $\{ u^i \in \{0, 1\}^3 \mid i = 1, \ldots, s \}$. For each $i$, $i = 1, \ldots, s$, let $R_i$ be the set of bits tested during $\text{comp}(u^i)$ before $v$ is reached. Let $R = \bigcup_{i=1}^s R_i$.

Proposition 4.8 For each bit $i$ outside $R$ and for each $j$, $j = 1, \ldots, s$, $w(u^j, v, C, P) = +$.

Proof: Let us develop a tree $T$ induced by $C$ from the node $v$. During $\text{comp}(u^i)$ after reaching $v$ the bit $i$ is either read for the first time - therefore in $T$ there is a branching node with the label $i$ and we have the desired $+$ on $i$, or the bit $i$ is not tested at all - we have a $+$ on $i$, too.

Let us remember the well-known fact that after $v$ no $u^i$ tests any bit from $R$. 11
Proposition 4.9 For each $i \in R$ if $i$ is a double-crossed bit for an $u^i \in C$ then $i$ is a double-crossed bit for all $w^i \in C$.

Proof: Let $i \in R$ be a double crossed bit for an $u \in C$. Then there is an $u' \in C$ such that $u'$ follows $u$ to its sink, $u, u'$ differ on $i$ and at least one from the inputs $u, u'$ is tested on $i$ before $v$.

Let us choose any $w^i \in C$. Let $w^i = u^i | \{1, 2, ..., n\} - R$. We see that $[u | R, w^i]$ and $[u' | R, w^i]$ have the same properties (mentioned above) as $u, u'$. Therefore they have a $\#$ on $i$. Moreover, they follow $w^i$ from $v$ to the sink in question. Therefore $w^i$ has a $\#$ on $i$, too.

\[ \square \]

Corollary 4.10 All $w^i \in C$ have the same region of the crossed (+), the double-crossed (#) and the non-crossed bits.

We see that $C$ is also an fd-class satisfying (R1),(R2).

\[ \square \]

Proposition 4.11 For each 1-branching program $P$ there is a branching program $P'$ such that

1. $f_{P'} = f_P$,
2. $|P'| \leq 3|P|$, 
3. $P'$ is a 1-branching program,
4. Each proper class of the distribution on $P'$ satisfies (R3) from the definition of the gentle branching programs.

Proof:

We will transform $P$ into $P'$ by a sequence of elementary steps. Each of them will consist in inserting one redundant test in the program. "Redundant test" means a test where both out-going edges go to the same node. It is clear that (1) $f_{P'} = f_P$ will be satisfied.

Let $Q$ be a 1-branching program arising from $P$ after some transformation steps. Let us describe the elementary step on $Q$.

For any input $u$ by $comp^P(u)$ we mean the initial part of $comp(u)$ from the source to the node at which $u$ is distributed. We say that an edge $e$ of $Q$ is supported by a path $p$ iff $e \in p$. Let $D(Q)$ be the subgraph of $Q$ given by edges which are supported by $comp^P(u)$, $u \in \{0, 1\}^n$. $D(Q)$ is also an oriented acyclic graph with one source. Its special properties are:

a) Its sinks are given by the nodes of $Q$ at which some inputs are distributed.

b) Except of sinks (and the source) all nodes of $D(Q)$ have in-degree 1 in $D(Q)$.

Let us consider the set of all nodes of $Q$ at which a proper class not satisfying (R3) is distributed. Let $v$ be a maximum node of this set (i.e. in $P$ $v$
has no predecessors from this set). Let \( F \) be the proper class distributed at \( v \). Let \( v_1, \ldots, v_l \) be the immediate predecessors of \( v \) in \( D(Q) \). In \( D(Q) \) for each \( j, 1 \leq j \leq l \) there is a unique path \( p_j \) going from the source to \( v \) via \( v_j \). Let \( M^j \) be the maximum original subcube going through \( p_j \) to \( v \). Let \( G^j \) be the set of bits from \( D \) (see def. 4.4) which are free for \( M^j \).

There is a \( j_0 \) such that \( G^{j_0} \neq \emptyset \) since \( F \) does not satisfy (R3). Let \( a \in G^{j_0} \).

On \( p_{j_0} \) let \( y \) be the first node such that the subprogram \( Q_y \) does not contain any test on \( a \). Such a \( y \) exists because in \( Q_v \) there is no test on \( a \) since \( a \in D \). \( y \) is not the source because in \( Q \) a test on \( a \) must exist since \( a \in D \). Let \( x \) be the immediate predecessors of \( y \) on \( p_{j_0} \). In \( Q_x \) there is a test on \( a \) therefore no path from the source to \( x \) contains a test on \( a \). Now it is clear that if we insert a redundant test (\( \geq \) two nodes) on \( a \) between \( x \) and \( y \) the resulting program \( Q' \) remains to be 1-branching program. At this moment we see that the point (3) of the statement of the Theorem is also satisfied.

Let us consider the distribution on \( Q' \). In comparison with the situation on \( Q \) there are changes concerning \( F \). In \( D(Q') \) the path \( p_{j_0} \) from \( D(Q) \) is interrupted at the second node of the test inserted between \( x \) and \( y \) where a new class \( F' \) of the distribution (on \( Q' \)) arises. It is clear that \( F' \) satisfies (R3) since inputs from \( F' \) have double-crosses (#) only on \( a \).

Slightly more complicated situation is at the node \( v \). If \( l \geq 3 \) than at \( v \) there is some class of distribution with \( l - 1 \) in-going original subcubes \( M^j \).

In the case \( l = 2 \) one path \( p_j \) with only one \( M^j \) remains alone. \( p_j \) expands to a tree whose branches end by meeting one another or by meeting another path in \( D(Q') \) or by reaching a sink.

We see that \( D(P') \) without the nodes with in-degree \( \geq 1 \) forms a tree. The first nodes of inserted redundant tests are placed only in the leaves of this tree. The remaining nodes are from \( P \). Hence, the number of the inserted first nodes is at most twice the maximum number of leaves of any tree in \( P \) which is equal to twice \((|P| + 1)/2 \). Since at least one leaf of our tree is from \( P \) we may estimate the number of inserted first nodes by \(|P| \). Hence, \(|P'| \leq 3|P| \) (the point 2 of the statement of the Theorem).

\(\square\)

From Proposition 4.7 it follows that \( P' \) is a gentle 1-branching program.

The proof of Theorem is completed.

\(\square\)
5 The lower bound

For the purposes of the present paper we put \( \alpha = \beta = 2^{2n^{1/8}} \) and \( \gamma = n^{1/8} \).

We shall define a Boolean function \( J \) for which we shall prove our lower bound for gentle b.p.'s.

Let us assume that \( n \) input bits are organized into a binary \((0, 1) (k\sqrt{n} \times \sqrt{n}/k)\)-matrix \( A \) where \( k \in N, k \geq 4 \) and moreover \( \sqrt{n}/k \) is an odd number. (Such a matrix \( A \) will be called "a matrix for \( J \).") The columns of \( A \) are indexed by numbers \( i \in <-(\sqrt{n}/k - 1)/2, (\sqrt{n}/k - 1)/2 > \). The column \( C_0 \) will be called the starting column.

In the following definition for two vectors (columns) \( M, C \) of the same dimension by \( M \oplus C \) (\( M \land C \), resp.) we shall mean the column which is the componentwise sum modulo 2 (componentwise conjunction, resp.) of \( M \) and \( C \). For a column \( M \) by \( |M|_1 \) we shall mean the number of 1's in \( M \). Similarly \( |M|_0 \).

On natural numbers larger than \( (\sqrt{n}/k - 1)/2 \) for \( q = 5n^{1/8} \) we define some intervals \( N_1, ..., N_q \) and a binary labeling of numbers from these intervals.

For \( i, 1 \leq i \leq q \), the labeling on \( N_i \) forms the word \((0^i1)^{12^i} \). We see that each two positions in \( N_i \) which differ by about \( i \) have different labels and that the summary length of \( N_1, ..., N_q \) is at most \( 24\gamma q^2 \). (The area of intervals \( N_i \) is adjacent to the matrix for \( J \).)

Further the interval \( k\sqrt{n}/4 - (\sqrt{n}/k - 1)/2; k\sqrt{n}/4 + (\sqrt{n}/k - 1)/2 > \) will be called the marked area \( MA \). We see that the inequality

\[ 24\gamma q^2 \leq k\sqrt{n}/4 - \sqrt{n}/k + 1 \]

holds. It gives more than a sufficient space for \( \bigcup_{i=1}^q N_i \) between the matrix for \( J \) and \( MA \).

Outside of \( \bigcup_{i=1}^q N_i \cup MA \) the odd positive numbers greater than \( (\sqrt{n}/k - 1)/2 \) will be labelled by 1, the even ones by 0. The negative numbers will be labelled in the symmetric way with only the exception that bits symmetric to those from \( MA \) will be labelled by zeroe's.

**Definition 5.1** Let \( A \) be a matrix for \( J \). Let \( s \) be an index of a column of \( A \), \( M \) be a column of dimension \( k\sqrt{n} \) and \( \square \in \{ \oplus, \land \} \).

We define \( \text{Jump}(M, s, \square) = (M', s', \square') \) as follows:

\[ M' =_{df} M \square C_s, \]

Let \( p =_{df} (|M'|_1 - |M'|_0)/2 \). If \( s + p \in <-(\sqrt{n}/k - 1)/2; (\sqrt{n}/k - 1)/2 > \) then \( s' = s + p \) and \( \square' = \oplus \).

If \( s + p \) is in \( MA \), then \( s' = s + p - k\sqrt{n}/4 \) and \( \square' = \land \).

If \( s + p \) is outside of the matrix for \( J \) and not in \( MA \), then \( s' = s + p \) and \( \square' \) remains undefined.

In the sequel \( M \) will be called the "memory column" incoming to \( C_s \) and \( M' \) will be called the "memory column" out-going from \( C_s \).
We see that if $s'$ is also an index of a column of $A$ we can iterate $\text{Jump}$ on $(M', s', \square')$.

**Definition 5.2** Let $A$ be a matrix for $J$. The value $J(A)$ is given as follows:

Starting with the values $M = \emptyset \sqrt{n}$, $s = 0$ and $\square = \emptyset$ we iterate $\text{Jump}$ until the number of iterations is $n$ or the last $s'$ is outside of $A(\cup MA)$. In the case of $n$ iterations we define $J(A) = 0$. In the other case $J(A)$ is defined as the label of $s'$.

It is easy to see that $J$ is computable on Turing machines within a polynomial time; $J \in \mathcal{P}$.

**Theorem 5.3** Each gentle branching program computing the function $J$ is of a size greater than $c(n) = 2^{2n/\alpha}$.

**Proof:** By contradiction. Let $P$ be a gentle program computing $J$, the size of $P$ is $c(n)$. Let us perform the distribution on $P$ and its factorization.

Since the number of sinks is at most $c(n)$ there are at most $c(n)$ improper classes of our distribution. Each improper class consists of only one original subcube otherwise there would be double-crossed bits and the class would be proper. Each original subcube $S$ in question has at most $2^k - (\sqrt{n}/k - 1)/2 + 1$ free bits on $C_0$ since in the other case there would be a possibility to reach both output values after the first jump which is in contradiction with the fact that all inputs from $S$ reach the same sink. Hence, $S$ has at least $b = c k \sqrt{n} - (\sqrt{n}/k - 1)/2 - 1$ fixed bits.

Therefore, the improper classes cover at most $c(n)2^n$ inputs. We see that there is a proper class of the distribution of the cardinality at least

$$(2^n - c(n)2^n) / c(n) \geq 2^{n/\alpha}.$$

Such a proper class is a significant one. Since $P$ is a gentle branching program at a node $v$, there is a significant proper class with a significant $fd$-subclass $F$ fulfilling the conditions $(R1), (R2), (R3)$ of the definition of gentle branching programs. The cardinality of $F$ is at least $f \geq 2^n/(\alpha, \beta)$.

Let $D_i, S_F, V_i, M_i, m_i, O_j$ be as in the Definition 4.4. We will prove that $D = \emptyset$; this will be a contradiction.

For each $i$ all inputs from $V_i$ have the same region $O_i$ of non-crossed bits and also the same region $R_i$ of bits with crosses +’s.

Let $a_i = d, |O_i|, d = d, |D|$, $r_i = d, |R_i|$. For $s$ an index of a column of the matrix for $J$ we define $O_{s,i} = d, O_i \cap C_s, D_s = d, D \cap C_s, R_{s,i} = d, R_i \cap C_s$ and further $a_{s,i} = d, |O_{s,i}|$, $d_s = d, |D_s|$, and $r_{s,i} = d, |R_{s,i}|$.

We assume w.l.o.g. that each node preceding $v$ in $P$ has out-degree $=1$. 

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Proposition 5.4 Let $a \in O^j_i$, $i \in \{1, \ldots, t\}$, $j \in \{1, \ldots, m_i\}$.

Then immediately before $v$ for all inputs from $M^j_i$ the bit $a$ is a non-crossed one.

More formally: Let $v_l$ be such a node immediately preceding $v$ that all inputs from $M^j_i$ go through $v_l$. Let $X_i$ be the set of all inputs from $F$ which reach $v_l$. Then for all $m \in M^j_i$ $(w(m, v_l, X_i, P))_a \in \{0, 1\}$.

**Proof:** By contradiction. There is an input from $M^j_i$ which has a) a cross + or b) a double-cross # on $a$ immediately before $v$.

Case a) On $a$ a cross may be placed if a test on $a$ is expected - but this is impossible since $a \in D$, or if during the computation from the source to the sink $a$ is not tested - but this is in contradiction with the fact that the original subcube $M^j_i$ is a maximum one.

Case b) If an $m \in M^j_i$ would have a double-cross # on $a$ before $v$, then $m$ is distributed to a node before $v$.

Let $\alpha^j_i$ be the cardinality of $O^j_i$. $O^j_i \cap C_s \alpha^j_i$ its cardinality.

Proposition 5.5 $\sum_{j=1}^{t} |M^j_i| (\alpha^j_i + \alpha_i)/f \leq log_2 c(n) + n - log_2 f$.

**Proof:** Let us consider our $f$ inputs from $F$. Immediately before $v$ they reach nodes $v_1, \ldots, v_z$ (each with out-degree = 1). These nodes define a partition of $F$ into classes $X_1, \ldots, X_z$.

According to the Theorem 3.6 we have

$(log_2 c(n)) \leq log_2 z \geq log_2 f - n + \sum_{j=1}^{t} (\sum_{x \in X_i} l w(x, v_l, X_i, P))/f$.

Therefore

$log_2 c(n) + n - log_2 f \geq \sum_{j=1}^{t} (\sum_{x \in X_i} \sum_{x \in M^j_i} l w(x, v_l, X_i, P)/f$,

$\geq \sum_{j=1}^{t} (\sum_{x \in X_i} \sum_{x \in M^j_i} (\alpha^j_i + \alpha_i))/f$ according to the previous proposition

$= \sum_{j=1}^{t} (|M^j_i| (\alpha^j_i + \alpha_i))/f$.

Corollary 5.6 $d \leq log_2 c(n) + n - log_2 f \leq 5n^{1/8}$.

**Proof:** (R3) implies that for all $i, j \alpha^j_i = d$. □

Corollary 5.7 $|W_1| \leq log_2 c(n) + n - log_2 f \leq 5n^{1/8}$.

Proposition 5.8 If $D_0 = \emptyset$, then $\forall s \neq 0 D_s = \emptyset$. □
Proof: By contradiction. Let us suppose \( \exists s \neq \emptyset \ D_s \neq \emptyset \). We are able to choose \( h_1, h_2 \) two assignments of \( D_s \) given by \( H_F \) such that or i) \( a = \lfloor (|h_1| - |h_1|/2 - (|h_2| - |h_2|)/2 > 0 \) or ii) there is a bit \( j \) ( a row \( j \) ) in \( D_s \) on which \( h_1, h_2 \) differ.

Case i).

We know from (R1) that \( |W_2| \leq \gamma + \log_2 e(n) + n - \log_2 f \leq 6n^{1/8} \) and from the Corollary 5.6 that \( |D| \leq \log_2 e(n) + n - \log_2 f \leq 5n^{1/8} \). Hence, there is an interval \( I \) of columns without any bit from \( W_2 \cup D \) containing \( 2(\gamma + \log_2 e(n) + n - \log_2 f) \leq 12n^{1/8} \) columns.

On \( S_F \) it is easy to see that for any specification of all bits outside of \( W_2(\cup D) \) there is a number of branches of \( S_F \) which respect this specification. Some possible branches may differ on some bits from \( W_2 \). If we choose such a branch \( b \) it is clear that in any case \( t \) induced by \( b \) is in \( T_F \) and therefore \( J([h_1, t]) = J([h_2, t]) \) (according to the condition (R2) of the definition of gentle). We shall construct a specification which induces only such \( t ^ \prime s \in T_F \) that on \( [h_1, t], [h_2, t] \) the jumps go from \( C_0 \) to \( I \) then to \( C_s \) then to the interval \( N_a \) and such that the positions reached by \( [h_1, t], [h_2, t] \) differ by \( a \). Hence, \( J([h_1, t]) \neq J([h_2, t]) \). This will be our contradiction.

Outside of \( I, C_0, C_s \), we choose an arbitrary assignment, say all zeroes.

For a set \( X \) of input bits by \( R(X) \) we mean the set of rows containing at least one bit from \( X \). For a set \( R \) of rows and a set \( C \) of columns by \( R \diamond C \) we mean the set of input bits which are in a row from \( R \) and in a column of \( C \).

We choose a set \( R_1 \) of rows such that \( R_1 \cap (R(W_2) \cup R(D)) = \emptyset \) and \( |R_1| = |I| \). We divide the rows except of \( R(W_2 \cup D) \) and \( R_1 \) into two sets \( R_2, R_3 \) such that \( (|R_2| + |R_1| - |R_3| - |R(W_2 \cup D)|)/2 = p_1 + k\sqrt{n}/4 \) where \( p_1 \) is the leftmost position of \( I \). We give one's into \( R_2 \diamond C_0, R_1 \diamond C_0 \), zero's into \( R_3 \diamond C_0, R(W_2 \cup D) \diamond C_0 - W_2 \). For any \( t \) induced by this specification it is clear that the first jump is to \( I \) (via \( MA \)) and that the next memory will be created by \( \land \).

On \( I \) first we give zeroes to \( I \diamond (R(W_2 \cup D)) \). We give one's to \( C, R_1, \) to \( C_{p_1+j} \diamond R_1 \) we give \( |R_1| - j \) one's and \( j \) zeroes (we have defined the content of \( I \diamond R_1 \)).

Further, we choose a set \( R_4 \) of rows such that \( R_4 \subseteq R_2, p_1 + (|R_3| + |R_2 - R_4| - |R_3| - |R_4| - |R(W_2 \cup D)|)/2 = s \). It is clear that such \( R_4 \) exists since \( R_2 \) is sufficiently large. We give ones to \( (R_2 - R_4) \diamond I \) and we give zeroes to \( R_4 \diamond I \). It is clear that for any \( t \) induced by this specification of bits of \( I \) the next jump is to \( C_s \). Explanation: Zero's in \( I \diamond R(W_2) \) eliminate the uncertainty about the contents of \( C_0 \diamond W_2 \). The specification of \( I \diamond R_1 \) eliminate the difference of positions of columns in \( I \). Moreover zeroes in \( I \diamond R(D) \) ensure that the memory column incoming to \( C_s \) has zeroes on \( R(D_s) \).
In any case the memory column $M$ incoming to $C_s$ has zeroes on $R(W_2 \cup D), R_4, R_3$ and $M$ has one’s on $R_2 - R_4$. The values of $M$ on $R_1$ depends on the case. The new out-going memory column $M'$ is created by $\oplus$.

First we give some values to $C_s \diamond R_1$, for instance zeroes. For any specification of $C_s \diamond (R_2 \cup R_3)$ it is clear that for any $t \in T_F$ induced by any branch of $S_F$ which respects this specification

a) There is an uncertainty concerning the jumped position of the size of at most $|R(W_2 \cup D) \cup R_1| < 4(\gamma + \log_2(n) + n - \log_2 f) \leq 24n^{1/8}$.

b) The positions for $[h_1, t], [h_2, t]$ differ about a since the incoming memory column $M$ has zeroes on $R(D)$.

So, we specify $C_s \diamond (R_2 \cup R_3)$ in such a way that $M'$ with supposed zeroes on $R(W_2 \cup D) \cup R_1$ would jump to the leftmost position of the interval $N_a$. In fact, in $R(W_2 \cup D) \cup R_1$ there are some ones in $M'$, but the jumped position still remains in $N_a$. Hence $[h_1, t], [h_2, t]$ have different values - contradiction.

Case ii)

We proceed in a similar way as in case i) till to the jump on $C_s$. We choose another sufficiently large interval $K$ of columns. From $C_s$ we jump via $MA$ to $K$. On $K$ we arrange zeroes and one’s in a similar way as in the previous case i) on $I$ with a substantial change that on the row $j \in R(D_s)$ on which $h_1, h_2$ differ we don’t give zeroes but ones. The effect of this change is that for any $t \in T_F$ induced by any branch of $S_F$ which respect our specification the positions jumped by $[h_1, t], [h_2, t]$ will differ by one. Therefore it suffices to arrange $K$ in such a way that in any case the next jump is to the interval $N_1$.

\[ \square \]

**Proposition 5.9 $D_0 = \emptyset$**

**Proof:** By contradiction. Let us suppose $D_0 \neq \emptyset$. We are able to choose $h_1, h_2$ two assignments of $D_0$ given by $H_F$ such that i) $a = \#f (|h_1| - |h_1|)/2 - (|h_2| - |h_2|)/2 > 0$ or ii) there is a bit $j$ (a row $j$) in $D_0$ on which $h_1, h_2$ differ.

Case i).

Outside of $W_2$ and $D$ we specify the bits of $C_0$ in such a way that with supposed zeroes in $W_2 \cap C_0$ and in $D_0$ the first jump would be to the leftmost position of the interval $N_a$. $N_a$ is sufficiently large, so for any content of $W_2 \cap C_0$ and $D_0$ the first jump must be in $N_a$. For any $t \in T_F$ induced by any branch of $S_F$ which respect our specification the positions jumped by $[h_1, t], [h_2, t]$ will differ about $a$. Hence, $J([h_1, t]) \neq J([h_2, t])$. A contradiction with (R2) of the definition of the gentle b.p.’s.

Case ii).

We proceed as in the case ii) of the previous proof. We find a sufficiently large interval $K$ of columns (outside of $W_2 \cup D$). Then we specify the bits in such a way that the jumps go from $C_0$ to $K$ via $MA$, from $K$ to $N_1$. 

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From the last two propositions it follows that \( D = \emptyset \). A contradiction. The theorem is completely proven.

\[\square\]

6 Upper bounds

6.1 The function \( f(A\vec{x}) \)

Let us suppose that our \( n \) input bits are organized into an \((l \times m)\)-matrix \( A \) and into an \((m \times 1)\)-vector \( \vec{x} \). We define \( f(A\vec{x}) = 1 \) iff \( A\vec{x} = \vec{0} \) (for all \( j, 1 \leq j \leq l \), \( \sum_{i=1}^{m} a_{ji}x_i = 0 \)).

For appropriate choices of \( l, m \), \( f \) is superpolynomially hard for \((1, +k)\)-branching programs [18]. This follows also from [10]. From [9] it follows that \( f \) is superpolynomially hard also for nondeterministic syntactic \( k \)-branching programs. It is clear that \( f \) is computable within polynomial time, \( f \in P \).

Theorem 6.1 There is a gentle branching program \( P \) which computes \( f \) and which is of the size \( O(n) \).

Proof:

Let us describe a natural branching program \( P \) which computes \( f \). \( P \) is a concatenation of branching programs \( P_1, \ldots, P_l \) where for each \( j = 1, \ldots, l \), \( P_j \) is responsible for the multiplication of the \( j \)-th row \( A_j \) of \( A \) with the vector \( \vec{x} \). Each \( P_j \) has two sinks; in one of them - with the label \( 0 \) - the computations on inputs for which \( A_j\vec{x} = \vec{1} \) stop, the other sink is stuck with the source of \( P_{j+1} \). Hence, \( P \) has \( l \) sinks with the label \( 0 \) and only one with the label \( 1 \). Now, let us describe \( P_j \). \( P_j \) is a levelled branching program of the width 2. \( P_j \) has \( m \) levels each with two tracks. In the upper tracks of \( l_i \) - there are two nodes with label \( x_i \) (with the only exception of the upper track of \( l_1 \) where only one node exists - the source of \( P_j \) labeled by \( x_1 \)). In the lower track of \( l_i \) there are two nodes labeled by \( a_{ji} \). On the last \( l_{m+1} \) level there are only two sinks.

In the upper track of each level \( l_i \), \( m \geq i \geq 2 \), one of the nodes with label \( x_i \) is called 0-node, the second one is called 1-node. In the upper track of the first level there is only the 0-node.

The 0-node in the upper track of the level \( i \) represents the inputs for which \( \sum_{s=1}^{m} a_{ji}x_s = 0 \) (i.e. it is reached by computations on these inputs). A similar rule for 1-node.

For \( a \in \{0, 1\} \) the 0-edge outgoing from the \( a \)-node on the level \( i \) reaches the \( a \)-node of the level \( i + 1 \). The 1-edge reaches one of the two nodes in the
lower track of the level $i$ labeled by $a_{ji}$. The edges out-going from each node in
the lower track reach the 0-node or 1-node of the level $i + 1$. They are arranged
in such a way that the rule concerning $a$-nodes is satisfied.

All $P_j$'s (and therefore $P$ too) are completely described. We see that $P$

computes $f$ and that $P$ is of the size of at most $|P_1| \leq 1.5m \leq 6n$.

Now, we shall prove that $P$ is a gentle branching program. Let us concen-
trate on $P_1$.

**Definition 6.2** $n_0^1 = 2^m$, for $i = 2, ..., m$, $a \in \{0, 1\}$, $n_a^i$ is the number of
inputs which reach the $a$-node of the level $i$. $n_a^{m+1}$ is the number of inputs
which leave $P_1$ by the sink stuckled with the source of $P_2$, $n_a^{m+1}$ for the other sink
of $P_1$.

**Proposition 6.3** $n_0^2 = 3n_1^2$.

*Proof:* The 0-node of the level 2 is reached by exactly those inputs $\overline{A\overline{E}}$ which
have $x_1 = 0$ or $x_1 = 1$ and $a_{11} = 0$. □

**Proposition 6.4** For $i = 2, ..., m$ if $n_a^i \geq n_a^i$, then $n_a^{i+1} \geq n_a^{i+1}$.

*Proof:*

\[ n_a^{i+1} = \frac{3}{4}n_a^i + \frac{1}{4}n_a^i \geq \frac{3}{4}n_a^i + \frac{1}{4}n_a^i = n_a^{i+1}. \]

The equalities follow from the fact that 0-edges from the $a$-node of the level
$i$ go directly to the $a$-node of the level $i + 1$ (for both $a = 0, 1$), and that 1-edge
outgoing from the $a$-node of the level $i$ goes to a node in the lower track of the
level $i$ where a branching on $a_{1i}$ is performed. □

**Corollary 6.5** $n_0^{m+1} \geq n_1^{m+1}$.

Let us perform the distribution on $P$. From our point of view the first three
nodes with in-degree > 1 will be interesting - the 0-node of level 2 $v_{02}$, the
0-node $v_{03}$ and the 1-node $v_{13}$ of the level 3. Let $M$ be the set of inputs such
that $\sum_{i=1}^m a_{1i}x_i = 0$. We know that $|M| \geq 2^m/2$.

**Proposition 6.6** At $v_{02}$ no $m \in M$ will be distributed there.

*Proof:* By contradiction. Le us suppose that there is an $m \in M$ which is
distributed at $v_{02}$. All $m'$ which follow $m$ from $v_{02}$ to a sink are in $M$ too.

For each $x_i, i = 1, ..., m$ in $v_{02}$ and after it all $m'$ following $m$ go through at
least one node with the test on $x_i$ ($x_1$ is tested at the source of $P_2$). Hence $m$
has no $\#$ on $x_1, ..., x_m$.

For each $a_{i,j}$ different from $a_{1j}$ after $v_{02}$ all $m'$ following $m$ simultaneously
test or don’t test $a_{i,j}$. Hence there is no $\#$ on $a_{i,j} \neq a_{1j}$. The only candidate
for \( \# \) is \( a_{1,1} \).

There are two possibilities .

a) \( m_{x_{1}} = 0 \)

Then for all \( m' \) which follow \( m \) from \( v_{02} \) to a sink \( m'_{x_{1}} = 0 \) holds due to the test on \( x_{1} \) at the source of \( P_{2} \). Hence for all such \( m' \) \( a_{1,1} \) is not tested at all, therefore there is + on \( a_{1,1} \).

b) \( m_{x_{1}} = 1 \)

Then for all \( m' \) which follow \( m \) from \( v_{02} \) to a sink \( m'_{x_{1}} = 1 \) holds (due to the test on \( x_{1} \) at the source of \( P_{2} \)). Since they all have \( x_{1} = 1 \) and they all reach \( v_{02} \), they all must have \( a_{1,1} = 0 \). Hence, \( a_{1,1} \) is noncrossed.

We see that there is no \( \# \) on \( m \) at \( v_{02} \) with respect to all inputs reaching \( v_{02} \). A contradiction.

\[ \square \]

**Proposition 6.7** At \( v_{02} \) there is a unique class of distribution containing all \( m \in \{0,1\}^{n} - M \) going through \( v_{02} \), the double-crossees \( \# \) are on bits \( x_{1}, a_{11} \).

**Proof:** If we develop the tree from \( v_{02} \) we see that all inputs, which in \( P_{1} \) go to the 0-sink, have double-crosses on \( x_{1}, a_{11} \). Therefore they are distributed at \( v_{02} \). \( \square \)

**Proposition 6.8** At \( v_{03} \) there is a unique fd-class (of the distribution) \( F_{0} \) containing only inputs from \( M \). These inputs have values \( x_{1} = 1, x_{2} = 1 \) and ((\( a_{11} = 0 \) and \( a_{12} = 0 \)) or (\( a_{11} = 1 \) and \( a_{12} = 1 \))).

**Proof:** All inputs from \( \{0,1\}^{n} - M \) which go through \( v_{03} \) and which are not distributed at \( v_{02} \) have the values \( x_{1} = 1, a_{11} = 1, x_{2} = 1, a_{12} = 1 \). Therefore at \( v_{03} \) there is no \( \# \) on these inputs.

Let us develop a tree \( T \) from \( v_{03} \). Since in \( P_{2} \) \( x_{1} \) and \( x_{2} \) are tested repeatedly, at any leaf of \( T \) the inputs reaching it have the same values of \( x_{1}, x_{2} \).

The case \( x_{1} = 0, x_{2} = 0 \). The bits \( a_{11}, a_{12} \) are not tested at all, therefore on \( a_{11}, a_{12} \) there are +’s (not \( \#'s \)). No double-crosses.

The case \( x_{1} = 0, x_{2} = 1 \). \( a_{11} \) is not tested at all. A cross + on \( a_{11} \). Since \( v_{03} \) is reached on \( a_{12} \) there must be the value \( a_{12} = 0 \). No double-crosses.

The case \( x_{1} = 1, x_{2} = 0 \). \( a_{12} \) is not tested at all - there is a cross on \( a_{12} \). Since \( v_{03} \) is reached, then there must be the value \( a_{12} = 0 \).

The case \( x_{1} = 1, x_{2} = 1 \). \( v_{03} \) is reached by inputs from \( M \) with the values \( x_{1} = 1, x_{2} = 1, a_{11} = 0, a_{12} = 0 \) or \( x_{1} = 1, x_{2} = 1, a_{11} = 1, a_{12} = 1 \). There are double-crosses (\( \# \)) on \( \{a_{11}, a_{12}\} \).
\[
\square
\]

**Proposition 6.9** At \(v_{13}\) there are two \(fd\)-classes. The first one contains inputs from \(\{0,1\}^n - M\), the double-crosses are on \(\{x_2, a_{12}\}\). The second one \(F_1\) contains inputs from \(M\), the double-crosses are on \(\{a_{11}, a_{12}\}\). These inputs have values \(x_1 = 1, x_2 = 1, a_{11} = 0, a_{12} = 1 \) or \(x_1 = 1, x_2 = 1, a_{11} = 1, a_{12} = 0\).

**Proof:** Let us develop a tree \(T\) from the node \(v_{13}\). We don’t consider the inputs from \(\{0,1\}^n - M\) which are distributed at \(v_{12}\). Each remaining input which goes to the 0-sink of \(P_1\) has the values \(x_1 = 1, a_{11} = 1\). Also we see that there are \#’s on \(x_2, a_{12}\).

Now let us consider the inputs from \(M\) reaching \(v_{13}\). Since in \(P_2\) the bits \(x_1\) and \(x_2\) are tested repeatedly, at any leaf of \(T\) the inputs from \(M\) reaching it have the same values of \(x_1, x_2\).

The case \(x_1 = 0, x_2 = 1\). \(a_{11}\) is not tested, on \(a_{11}\) there is a cross. \(a_{12}\) must be equal to 1. No double-crosses.

The case \(x_1 = 1, x_2 = 0\). \(a_{12}\) is not tested, \(a +\) on \(a_{12}\). \(a_{11}\) must be equal to 1. No double-crosses.

The case \(x_1 = 1, x_2 = 1\). For reaching \(v_{13}\) it must hold \(a_{11} = 0, a_{12} = 1\) or \(a_{11} = 1, a_{12} = 0\). There are double-crosses on \(\{a_{11}, a_{12}\}\).

\[
\square
\]

Let \(F\) be the largest class from \(F_0, F_1\).

**Proposition 6.10** The cardinality of \(F\) is at least \(2^n/20\).

**Proof:**

We divide all inputs into the groups \(G_1, ..., G_{16}\) according to the equalities \(x_1 = c_1, a_{11} = c_2, x_2 = c_3, a_{12} = c_4, c_i \in \{0,1\}\). Ten of them \(G_1, ..., G_{10}\) reach \(v_{13}\), \(G_{11}, ..., G_{16}\) reach \(v_{13}\). For each \(i, i' = 1, ..., 10\) we define \(G_i^0 = G_i \cap M\). For \(i = 1, ..., 10\) all \(G_i^0\) have the same cardinality since in and after \(v_{13}\) the computation depends only on \(x_3, x_{13}, a_{13}, ..., a_{1m}\). Similarly for \(G_{11}^0, ..., G_{16}^0\).

Let \(M_0\) \((M_1, \text{resp.})\) be the set of inputs from \(M\) which go through \(v_{13}\) \((v_{13}\), resp.\). \(|M_0| + |M_1| = |M| \geq 2^n/2\) (see Corollary 6.5).

At \(v_{13}\) the inputs from two groups from \(G_1^0, ..., G_{10}^0\) form \(F_0\). At \(v_{13}\) the inputs from two groups from \(G_{11}^0, ..., G_{16}^0\) form \(F_1\).

Hence, \(|F_0| = |M_0|/5\), \(|F_1| = |M_1|/3\). If \(|M_0| \geq |M_1|\) then \(|F_0| = |M_0|/5 \geq |M|/10 \geq 2^n/20\). If \(|M_1| \geq |M_0|\) then \(|F_1| = |M_1|/3 \geq |M|/6 \geq 2^n/12\). Hence, \(|F| > 2^n/20\).

\[
\square
\]

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Corollary 6.11 At least one of the distribution classes at $v_{03}$ and $v_{13}$ is a significant one.

Now it suffices to prove that $F$ is a significant $fd$-class which satisfies (R1),(R2),(R3) from the definition of the gentle programs.

First, let us consider the case $F = F_0$. We see that $F_0$ is a significant subclass of $F_0$.

ad (R1). Let us develop the tree $S_F$ from $v_{03}$. We see that $x_1 = 1, x_2 = 1$ are non-closed bits on each branch of $S_F$. The other bits which may sometimes be non-closed are $a_{1,m}, x_m$.

ad (R2). (R2) is clearly satisfied since all $t \in T_F$ have $x_1 = 1, x_2 = 1$. Hence, the result does not depend on $a_{1,1} = 0, a_{1,2} = 0$ or $a_{1,1} = 1, a_{1,2} = 1$.

ad (R3).

Let $D$ be the set of double-crossed bits. We know that all branches of $S_F$ have $D = \{a_{1,1}, a_{1,2}\}$. From Proposition 6.7 it is clear that in $D$ there are no original subcubes. Hence, for all $i, j O_i^j = D$.

Similarly for the case $F = F_1$.

The theorem is proven.

\[ \Box \]

6.2 Okolnishnikova function $F_{n,s}$

Let us describe the functions $F_{n,s}$, $n, s$ natural numbers, $n > s$, $s$ divides $n$, which are the witness functions for the hierarchy of syntactic $k$-branching programs with respect to $k$ ([15],[16]).

There are $N = \binom{n}{s}$ variables, each of them is indexed by a subset $a$, $|a| = s$, of the set $I_n = \{1, \ldots, n\}$. For $i = 1 \ldots n$, $W_i = \{a | a \subseteq I_n, |a| = s\}$. We see that $|W_i| = \binom{n-1}{s-1}$. We define $F_{n,s} = \bigwedge_{i=1}^{N} \bigvee_{a \in W_i} x_a$.

Theorem 6.12 There is a gentle program $P$ (on $N = \binom{n}{s}$ variables) such that a) $P$ computes $F_{n,s}$, b) $|P| = O(N \log N)$.

Proof:

We shall describe a branching program $P$ of the size $N \log N$ which computes $F_{n,s}$ such that $P$ is a gentle branching program.

Let $P_i$ be a program of a standard form which computes the function $\bigvee_{a \in W_i} x_a$. $P_i$ has two sinks; in the source of $P_i$ there is a test on the first $x_a$, the 1-edge of this test ends in the 1-sink, the 0-edge leads to the test of the next variable (or in the case of the test on the last variable to the 0-sink). On its turn the next variable will be managed in the same way. We see that $|P_i| = \binom{n-1}{s-1} + 2$.

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Roughly speaking, $P$ will be a concatenation (a conjunction) of $P_i$’s with some additional properties.

We choose a variable $x$. Let all $P_i$’s testing $x$ have the tests on $x$ at their sources. Moreover all $P_j$’s with tests on $x$ will form the initial sequence of our concatenation of all $P_i$’s. (The number of $P_i$’s with a test on $x$ is $s$.)

From technical reasons after this initial sequence we insert a chain of redundant tests of all variables different from $x$ into the concatenation of all $P_i$’s. The tests of this chain are of specific form so that both 0-edge and 1-edge of a test of a variable go to the same node where the next variable is tested (in the case of the last test they go to the source of the first subprogram from the remaining $P_i$’s).

Later we will define a linear ordering on variables tested in the initial sequence of $P_i$’s.

Now we see that $P$ computes $F_{n,s}$ and that $|P| \leq n|P_i| + N \leq O(N \log N)$.

Let us investigate the distribution on $P$. Let $v_1, ..., v_s$ be the sinks of $P_i$’s from the initial sequence such that $v_i$’s are sources of the next parts of $P$. $v_1, ..., v_{s-1}$ are the unique candidates for a class of distribution before $v_s$ (since they are the unique nodes before $v_s$ with indegree $> 1$). If an input of $m \in \{0, 1\}^N$ belongs to a class of the distribution at a node $v_1, ..., v_{s-1}$, then $m$ does not have $\#$ on $x$ since in $v_1, ..., v_{s-1}$ a test on $x$ follows. Moreover, $\comp(m)$ does not reach $v_s$ since the tests on all variables different from $x$ follow. Hence, each $m$ which is distributed at $v_1, ..., v_{s-1}$ ends in a 0-sink of one of $P_1, ..., P_s$. The number of such $m$’s is not greater than $s 2^{N/2^{s-1}}$.

The class of the distribution at $v_s$ is a larger one. Each input $m$ with $m_x = 0$ which reaches $v_s$ has a double-cross $\#$ on $x$ since $v_s$ is also reached by the input $m'$ which differ from $m$ only on $x$. Further, at $v_s$ the double-cross $\#$ may be only on $x$ - this follows from the fact that in the next chain there are tests on all variables different from $x$.

We see that the class of the distribution at $v_s$ is a significant proper class which consists from the unique $fd$-subclass $F$.

It is clear that $F$ satisfies (R2) from the definition of gentle programs. It remains to be verified that (R1),(R3) are satisfied, too.

First we shall complete the specification of the construction of $P$. We define a linear ordering $\prec$ on the set of all variables (different from $x$). The initial part of $\prec$ is given by the (arbitrary) ordering of tests in $P_1$. In $P_2$, first the tests on variables which are tested also in $P_1$ respect $\prec$. Moreover they precede the tests on the remaining variables. In $\prec$, these remaining variables follow after the previous ones and they respect the test ordering in $P_2$. In each $P_i$, the variables from the previous $P_j$’s are tested in respect to $\prec$ at first, then new
variables are tested and $\prec$ is extended.

Let $y_1, \ldots, y_s$ be the variables which in $P_1, \ldots, P_s$ are tested as the last variables. Let us suppose that we have $m \in F$. For each variable $y \neq x$ we shall find out whether $y$ is crossed ($\pm$) at $v_s$ or not.

If $y$ is not tested during $\text{comp}(m)$ between the source of $P_1$ and the node $v_s$, then $y$ is crossed since there is $m', m'_y \neq m_y$ and $m'_y = m'_{\neq y}$. Both $\text{comp}(m), \text{comp}(m')$ reach $v_s$ and they must branch on $y$ in the chain after $v_s$.

Now, we have $m \in F$, $y \neq x$, and $y$ is tested during $\text{comp}(m)$ between the source of $P_1$ and $v_s$.

We have two cases A) $m_x = 0$, B) $m_x = 1$.

ad B). Since $m_x = 1$ and $m \in F$ there is an $m'$ such that $m'_x = 0$ and $m'$ reaches $v_s$. From $v_s \text{comp}(m)$ and $\text{comp}(m')$ follows the same path, hence they have the same set of crosses $\pm$’s. We see that B) is the same as A) for $m'$.

ad A). We have $m \in F$, $m_x = 0$. Let $y \neq x$ be a variable tested during $\text{comp}(m)$ before $v_s$.

a) the case $m_y = 0$. There is an input $m'$ such that $m'_y = 1$ and $m'_{\neq y} = m'_{\neq y}$.

We see that $m'_x = 0$ and $m'$ reaches $v_s$ ($m'$ has more one’s than $m$). Therefore, $m' \in F$ (see above). $m$ and $m'$ branch on $y$ in the chain after $v_s$. Hence, $m, m'$ have $+$ on $y$.

b) the case $m_y = 1$. There is an input $m'$ such that $m'_x = m'_x$, and $m'_y = 0$ and $m'_{\neq y}$ is an unary word over $\{1\}$. If $y \neq y_1, \ldots, y_s$ (last variables), then $m'$ reaches $v_s$. Since $m'_x = 0$ we have $m' \in F$. $m$ and $m'$ branch on $y$ in the chain after $v_s$. Hence, both $m, m'$ have $+$ on $y$.

We see that for all input $m \in F$ all non-crossed bits are from the set $\{y_1, \ldots, y_s\}$. Since $s < N^{1/8}$, (R1) is satisfied.

(R3) is satisfied trivially since $D = \{x\}$ and in any computation on any input $x$ is tested before $v_s$.

\[\square\]

### 6.3 Multipointer functions

In [20] it was proven for $k(n)$ up to $k(n) = 2^{-1}, n^{1/8}, \log^{-1/4}n$ that multipointer functions $f_{n,k(n)}$ causes the hierarchy for $(1, +k(n))$-branching programs - it means that $f_{n,k(n)}$ is polynomially easy for $(1, +k(n))$-branching programs but superpolynomially hard for $(1, +k(n) - 1)$-branching programs. We shall prove that $f_{n,k(n)}$ (for $k(n)$ till $n^{1/8}$) are polynomially easy for the gentle branching programs.
We define the functions $f_{n,k}$ as follows: First we describe them informally. The $n$ variables are divided into $k$ blocks of length $m$. For every $j = 1, 2, \ldots, k$, a weighted sum of the bits of block $j$ determines an index $i_j$ of some of the input bits. Then the value of the function is the parity of the bits determined by $i_j$ for $j = 1, 2, \ldots, k$. The exact definition of $f_{n,k}$ requires some technical notation.

For every natural number $n$, let $p(n)$ be the smallest prime greater than $n$. Consider the set $\{1, 2, \ldots, n\}$ as a subset of $Z_{p(n)}$, the field of the residue classes modulo $p(n)$. Then for every $t \in Z_{p(n)}$ let $\omega(t) = t$ if $t \in \{1, 2, \ldots, n\}$ and $\omega(t) = 1$ otherwise.

**Definition 6.13** For every $t = (t_1, \ldots, t_k) \in \{1, \ldots, n\}^k$ and every $x \in \{0, 1\}^n$, let $\text{Par}(x,t) = \oplus x_{t_1} \oplus \ldots \oplus x_{t_k}$.

**Definition 6.14** Let $k$ divide $n$ and let $m = n/k$. Let $\psi_{n,k} : \{0, 1\}^n \rightarrow \{1, 2, \ldots, n\}^k$ be defined as follows. For every $x \in \{0, 1\}^n$ let $\psi_{n,k}(x) = (t_1, \ldots, t_k)$ where for every $j = 1, 2, \ldots, k$ $t_j = \omega(\sum_{i=1}^{m} i x_{(j-1)m+i})$ where the sum is evaluated in $Z_{p(n)}$. Moreover let $f_{n,k}(x) = \oplus \text{Par}(x, \psi_{n,k}(x))$.

**Theorem 6.15** For $k < n^{1/8}$ $f_{n,k}$ is computable on a gentle branching program of the size $O(n^2)$.

**Proof:**

Given $k \in N$ we shall construct a branching program $P$ computing $f_{n,k}$. The structure of $P$ will reflect $k$ blocks of input bits with the corresponding pointers.

$P$ will consist from $k$-parts $P_i$, $i = 1, \ldots, k$. Each $P_i$ will have two input nodes (with the exception of $P_1$ which has only one input node - the source of $P$) and two output nodes. The output nodes of $P_k$ are the sinks of $P$. The two output nodes of $P_i$ are stucked with the input nodes of $P_{i+1}$ for all $i = 1, \ldots, k-1$. One of the output node of $P_i$ is called the 0-node, the other is called the 1-node. For $a \in \{0, 1\}$ it is the truth that the $a$-node of $P_i$ is reached by all inputs for which the value of the parity of the first $i$ bits pointed out by pointers from the first $i$ blocks is equal to $a$.

$P_1$ has $m + 1$ levels. On each level from the first $m$ ones one bit is tested. Each output node of the $i$-th level represents one magnitude of the contribution of the first $i$ bits to the pointer. On the last level ($m + 1$) the bits which are pointed out are tested. The outgoing edges are stucked to the $a$-nodes of $P_1$.

The other $P_i$'s consist from two copies of $P_1$ (labelled by bits from the $i$-th block). The only change is on the $m+1$-st levels where the outgoing edges must point to the correct $a$-node.

We see that that $P$ computes $f_{n,k}$ and that the size of $P$ is not greater then $(2k-1)|P_1| \leq O(n^2)$.
It remains to prove that $P$ is a gentle branching program. We assume that at the top of $P_1$ there is a small irregularity. The initial part of $P_1$ consists of the full tree of depth 3 where at the root of the tree the variable $x_2$ is tested, at its successors $x_3$ is tested and the last tests are on $x_5$. The leaves of the branches $(x_2 = 1, x_3 = 1, x_5 = 0)$ and $(x_2 = 0, x_3 = 0, x_5 = 1)$ are sticked into one node $v$ (they give the same contribution to the pointer on the first block). Moreover we suppose that at the level $m$ of $P_1$ the variable $x_m$ is tested.

Now let us perform the distribution on $P$. The node $v$ is one of the first candidates for nodes at which some clases are distributed. Let us consider an input $p$ going through $v$ and assume that $p_2 = 1, p_3 = 1, p_5 = 0$. Let us take the input $p'$ with $p'_2 = 0, p'_3 = 0, p'_5 = 1$ such that $p'$ follows $p$ to a sink. We see that $p, p'$ have “s” on $x_2, x_3, x_5$. Since during the path through $P$ all variables different from $x_2, x_3, x_5$ are tested $p, p'$ must equal on these variables. Moreover $p, p'$ must avoid possible tests on $x_2, x_3, x_5$ at the last levels of $P_1$’s. Let us estimate the number of such $p$ - i.e. the cardinality of the class $F$ of the distribution at $v$.

We see that from the pair of inputs which differ only in the last variable of the first block at most one of them goes through a test on one from the variables $x_2, x_3, x_5$ at the last level of $P_1$ since the possible contribution of the last variable to the pointer is very large. The same fact holds for all remaining blocks. Hence the cardinality of $F$ is at least $\left(2^n/4\right)2^{-1}2^{k-1} \geq 2^n/2^{2n^{1/8}}$. Therefore $F$ is a significant class of the distribution.

We see that $F$ is also an $fd$-class. It remains to verify the conditions (R1),(R2),(R3) from the definition of the gentle branching programs.

Ad (R1). In each $P_i$ we want to avoid the tests on $x_2, x_3, x_5$. This can be ensured by fixing the last variable in each block. Hence at most these $n^{1/8}$ variables are noncrossed in any computation on input from $F$.

Ad (R2). Since after $v$ the inputs from $F$ are not tested on $x_2, x_3, x_5$ the reached sink does not depend on these variables.

Ad (R3). Clearly.

$$\square$$

7 Other upper bounds

7.1 Parity of triangles in the graph

Let $G = (V, E)$ be a graph, $V = \{v_i\}_{i=1}^m$, $E \subseteq V \times V$. $G$ defines a binary matrix $A = (a_{ij})_{i,j=1}^m$ where $a_{ij} = 1$ iff $(v_i, v_j) \in E$. 

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Since $G$ is unoriented, $A$ is symmetric. By a code of $G$ we mean the binary string $e = a_{12}...a_{1m}a_{22}...a_{2m}...a_{j,j+1}...a_{jm}...a_{m-1,m}$.

Let $f$ be a Boolean function such that $f(e) = 1$ if and only if $e$ is a code of a graph which has an odd number of triangles.

For $f$ the lower bounds $2^{n^3/2}$ on the size of 1-branching programs have been proven in [2],[22].

**Theorem 7.1** There is a gentle branching program $P$ of size $O(n^{3/2})$ which computes $f$.

**Proof:** In $P$ the following algorithm is implemented.

We consider all possible triangles (on $m$ vertices), each (possible) triangle is considered only once.

We fix an (possible) edge, say $e_1$, arbitrarily. In the first phase of the algorithm we consider all (possible) triangles over $e_1$. In the second phase, we consider the remaining triangles.

$P$ will have $k$ levels $P_i$, $i = 1..k$ where $k$ is the number of all possible triangles. Each level $P_i$ will be responsible for considering a triangle.

Each $P_i$ has two inputs, two outputs and some internal nodes (with exception of $P_1$ with one input node - the source of $P$). The output nodes of $P_{i-1}$ will be stuck with the input nodes of $P_i$. The output nodes of $P_k$ are sinks of $P$.

For $i = 1..k$, one output node of $P_i$ is called 0-node, the other one the 1-node. The following rule is fulfilled: Starting at the source of $P$, each input which among the first $i$ considered possible triangles has an even (odd, resp.) number of actual triangles achieves the 0-node (the 1-node, resp.) of $P_i$.

We see that $P$ computes $f$. It remains to specify the internal nodes of each $P_i$. For each $i$ the input nodes of $P_i$ will be the roots of two copies of the full balanced tree of depth 3 with tests on variables which represent the edges of which the $i$-th triangle consists. The out-going edges in the depth 3 point to 0- and 1-node of $P_i$ in such a way that the rule concerning the 0- and 1-nodes is satisfied.

We see that the size of $P$ is at most $O(n^{3/2})$.

Now we are going to prove that $P$ is a gentle b.p.. Let $i_0$ be the last level of $P$ where the edge $e_1$ is tested. Let us perform the distribution on $P$ and let us consider the situation at the 0-node and 1-node of $P_{i_0}$. In $P$ only the output nodes of $P_i$’s have in-degree $> 1$, hence only at the output nodes the double-crosses ($\#$) may arise. From the construction of $P$ it is clear that in $P_{i_0+1},...,P_k$
all computations have tests on all variables except $e_1$. Hence in $P_1, ..., P_{i_0}$ the 
$\#$'s may be at most on $e_1$. Moreover, in $P_1, ..., P_{i_0-1}$ the $\#$ on $e_1$ is impossible 
since a test on $e_1$ follows in $P_{i_0}$.

Let us investigate the possibility of $\#$ on $e_1$ at the output nodes of $P_{i_0}$. 
Let $m$ be an input. $m$ has $\#$ on $m_1$ iff there is $m'$ such that a) $m_1' \neq m_1$, b) 
$\text{comp}(m')$ follows $\text{comp}(m)$ from the node in question to a sink. In the output 
nodes of $P_{i_0}$ the condition b) implies that $m_{\neq 1} = m'_{\neq 1}$ since in the remaining 
part of $P$ it tests all variables different from $e_1$.

W.l.o.g. we suppose that $m_1 = 1, m'_1 = 0$. We see that $m$ is the code of 
such a graph $G_m$ that deleting the edge $e_1$ from $G_m$ the parity of the number of triangles does not change since $\text{comp}(m')$ goes to the same sink as $\text{comp}(m)$. 
Hence, in $G_m$, the number of triangles with the edge $e_1$ is even. Hence, $\#$ may 
only be at the 0-node of $P_{i_0}$.

On the other hand, the condition "$m_1 = 1$ and in $G_m$ there is an even 
number of triangles with $e_1$" is a sufficient one for a $\#$ on $m_1$.

Let $F$ be the set of all inputs which have a $\#$ on $e_1$ at the 0-node of $P_{i_0}$. 
Later we shall see that $F$ is a significant class of the distribution. It is clear that $F$ is also an $\text{fd}$-class and that for $F$ the condition (R2) from the definition 
of gentle b.p.'s is satisfied.

For verifying (R1), from the 0-node of $P_{i_0}$ let us develop the tree $S_F$ induced 
by $F$.

Since during each computation each variable different from $e_1$ is tested there 
are exactly two inputs $m, m', m_1 \neq m'_1, m_{\neq 1} = m'_{\neq 1}$ in each leaf of $S_F$.

Let us arbitrarily choose a branch $b$ of $S_F$. We shall investigate the set of 
non-crossed bits of inputs following $b$.

What concerns variables which have not been tested in $P_1, ..., P_{i_0}$ there will 
be a cross on each of them. It is clear since after $P_{i_0}$ each computation goes 
through nodes with tests on all (till this moment) non-tested variables.

Let $m \in F$ be an input following $b$, let $x \neq e_1$ be a variable tested in 
$P_1, ..., P_{i_0-1}$. Let $y$ be such a variable that the edges represented by $e_1, x, y$ form 
a triangle. W.l.o.g. we suppose $m_1 = 1$.

If $m_y = 0$ we define an input $m'$ as follows : $m'_{\neq x} \overset{=} \neq m_{\neq x}$ and $m'_x \neq m_x$. 
Since $y = 0$ it is clear that $m'$ has the same number of triangles over $e_1$ as $m$ 
does (even) and therefore $m' \in F$. $m'$ will branch with $m$ on $x$ therefore $m$ has 
a cross on $x$.

If $m_y = 1$, then we choose $m'$ such that 
(1) $m'_{e_1} \overset{=} \neq 1$, 

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(ii) \( m'_{x} \neq m_{x} \),
(iii) \( m' \) equals \( m \) on variables which are tested on \( b \) before \( x \),
(iv) On the other variables with exception of those which represent the \( i_{0} \)-th triangle over \( e_{1} \) we define \( m' \) arbitrarily.
(v) On the remaining two variables we define \( m' \) in such a way that \( m' \) has an even number of triangles over \( e_{1} \).

Now we see that \( m' \in F \), \( m' \) follows \( m \) along \( b \) till to the test on \( x \). Hence \( m \) has a + on \( x \).

We see that on each branch \( b \) of \( S_{F} \) the only candidates for non-crossed bits are the two bits which are tested at the level \( i_{0} \). Hence \((R1)\) is satisfied.

By the way we have obtained the fact that the cardinality of \( F \) is at least \( 2^{n-3} \). Hence \( F \) is a significant class of the distribution on \( P \).

\((R3)\) is satisfied clearly.

\( \square \)

### 7.2 Half-cliques-only

By a half-clique-only we understand an unoriented graph where one half of vertices form the full graph and the other vertices are isolated, \((G = (V, E), V = V_{1} \cup V_{2}, |V_{1}| = |V_{2}|, E = V_{1} \times V_{1})\).

In 1984 [24] it was proven that the Boolean function \( f \) which gives 1 exactly on the codes of half-cliques-only is subexponentially hard for \( 1 \)-branching programs. In 1993 [4] it was proven that \( f \) is subexponentially hard even for nondeterministic \( 1 \)-branching programs.

By a code of a graph we mean the same string \( e = s_{1}, ..., s_{m-1} \) as in the previous subsection where \( s_{j} = a_{i_{j}, j+1}, ..., a_{i_{j}, m} \).

Let us notice that the code of any half-clique-only has some special properties:

a) Exactly \( m/2 - 1 \) segments \( s_{i} \), \( i = 1, ..., m \), contain at least one symbol 1,

b) For \( i < j \) if \( a_{i, j} = 1 \) then the segment \( s_{j} \) is an suffix of the segment \( s_{i} \).

Let us introduce a natural algorithm computing \( f \).
1. Search for the first segment \( s_{i} \) containing a symbol 1; \( s_{i} \) is pointed out.
2. Move 1 to a counter \( C \).
3. Perform \( Proc \) until \( C = m/2 - 1 \).
4. Let \( s_{j} \) be the last segment which is pointed out. Check whether the segment \( s_{k} \) is a zero-sequence for each \( k, k > j \).

\( Proc: \)
a) Let $s$ be the segment which is pointed out. Search for the first $j_0, j_0 > j$
such that $a_{j_j} = 1$ (the leading one).
b) Check whether $s_k$ is a zero sequence for each $k, j < k$.
c) Check whether $s_{j_j}$ is a suffix of $s_j$.
d) $s_{j_j}$ is pointed out.
e) Add 1 to $C$.

It is easy to see that such an algorithm can be implemented in a branching
program $P_1$ of the size of at most $O(m^4) \leq O(n^2)$.

**Theorem 7.2** $f$ is computable on a gentle branching program $P$ of the size
$O(n^2)$.

**Proof:**
At the top of $P$ there is a tree of the depth 3 with branching on $a_{1,2}, a_{1,3}, a_{2,3}$.
The branches $a_{1,2} = 1 a_{1,3} = 0 a_{2,3} = 1$ and $a_{1,2} = 1 a_{1,3} = 1 a_{2,3} = 0$ are stucked
into one node $v$ which is a sink of $P$ with the label 0. The other branches are
stucked to a node which is the source of $P_1$.

It is clear that $P$ computes $f$ and that $|P| \leq O(n^2)$. It remains to prove
that $P$ is a gentle program.

The sink $v$ is reached by $2^n/4$ inputs. Each of them has \#’s on $a_{1,3}, a_{2,3}, 1$
on $a_{1,2}$ and +’s on the other bits. Hence, these inputs form a significant proper
class of the distribution. It easy to see that this class is an $fd$-class fullfilling
the conditions $(R1),(R2),(R3)$ from the definition of gentle branching programs.

$\square$

### 7.3 Ablayev-Karpinski’s function

In [1] there is defined a function $f$ on words over four-letter alphabet \{0, 1, 0, 1\}
as follows : $f(u) = 1$ iff $p(u) = p(u)$ where $p, \bar{p}$ are projections such that
$p(u_1u_2) = p(u_1)p(u_2)$ and $p(0) = 0, p(1) = 1$ and $p(\bar{0}) = p(\bar{1}) = 0$ (= the empty
word). Similarly for $\bar{p}$.

In [1] the authors prove that $f$ is superpolynomially hard for nondeterministic
four-letter 1-branching programs. We use a variant of $f$ where the four
letter alphabet is encoded by strings of the length two \{00, 10, 01, 11\} where
the first bit represents the value and the second one represents the type.

There is a natural algorithm $A$ computing $f$ which is based on the procedure:
"Search for the next value of the first type and search for the next value
of the second type, and compare".

Such an algorithm can be implemented in a branching program $P_1$ of the
size of at most $O(n^3)$.

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Theorem 7.3 $f$ is computable on a gentle branching program $P$ of size $O(n^3)$.

Proof: At the top of $P$ there is a tree on depth four which branches on the first four variables. The branches induced by inputs with prefixes 0011, 1001, 0110, 1100 are stucked into a node $v$ which is a sink of $P$ with the value 0. The other branches are stucked into the other node which is the source of $P_1$.

It is clear that $P$ computes $f$ and that the size of $P$ is at most $O(n^2)$. It remains to prove that $P$ is a gentle program.

The sink $v$ is reached by $2^n/4$ inputs. At $v$ each of them has #’s on the first four bits. The other bits of these inputs have +’s. We see that the class of these inputs is a significant proper class of the distribution, moreover it is also a significant fd-class fulfilling the conditions (R1),(R2),(R3) from the definition of gentle branching programs.

\[\Box\]

7.4 Multiplication

S. Ponzo [17] has proven a $2^{\sqrt[3]{n}}$ lower bound for 1-branching programs computing a function $f$ which is closely connected with the multiplication. For the same function $f$ we shall prove a $O(n^3)$ upper bound for gentle branching programs.

Given $x, y \in \{0, 1\}^m$, $x = x^m x^{m-1}...x^2 x^1$, $y = y^m y^{m-1}...y^2 y^1$ we define $f(xy) = y^m z^m$ where $z^m$ is the $m$-th digit (from the right) of the binary number $z = x \times y \in \{0, 1\}^{2m}$, $z = z^{2m} z^{2m-1}...z^2 z^1$.

Theorem 7.4 $f$ is computable on a gentle branching program $P$ of the size $O(n^3)$.

Proof: In $P$ we will implement the well-known algorithm for the multiplication. The algorithm sums the rows $R_i = r_{i+1}...r_{i+1}$ for $i = 1...m$ where $r_{ij} = df x^j \times y^i$. We will consider the columns $C_j = c_{1,j}...c_{m,j}$ where $c_{ij} = df r_{i,j-1}+1$ for $j - i + 1 > 0$ and $c_{ij} = 0$ otherwise. The value $f(xy)$ is given as the sum of the last digit of the sum of $C_m$ and of the last digit of the transition from the previous columns $C_m-1,...,C_1$.

The structure of $P$ will reflect the columns $C_1,...,C_m$ and the magnitude of the respective transitions.

$P$ will consists of the parts $P_1,...,P_m$. With exception of $P_m$, $P_j$’s are of the regular form as follows:

Each $P_j$ has some input nodes ad some output nodes. Each input node represents one magnitude of the transition after the sum of the columns $C_1,...,C_{j-1}$.
It means that all inputs $xy$ having the transition of the same magnitude (after $C_1, ..., C_{j-1}$) go through the same input node of $P_j$. We see that $P_1$ has only one input node - the source of $P$. The output nodes of $P_j$ are stucked with the input nodes of $P_{j+1}$.

Each $P_j$, $j = 1, ..., m-1$ consists of $l_j$ levels. Each level $L_i$, $i = 1, ..., l_j$ corresponds to the value $c_{ij}$. Each level $L_i$, $i = 1, ..., l_j$ has some input nodes and some output nodes. Each input node represents one magnitude of the sum of the transition after the columns $C_1, ..., C_{j-1}$ and of $c_{1j}, ..., c_{i-1j}$. The output nodes of $L_i$ are input nodes of $L_{i+1}$.

At the level $L_i$ it is necessary to find out the value $c_{ij}$. Each input node of $L_i$ is the root of the full tree of the depth 2. The tree tests the variable $x^{j-i+1}$ at its root and the variable $y^i$ at the two immediate successors of the root. The leaves of these trees are stucked into the output nodes of $L_i$.

Now $P$ is described till $P_{m-1}$. Since the maximal magnitude of the transition is not larger than $2m$ the inequality $\sum_{j=1}^{m-1} |P_j| \leq O(n^3)$ holds.

It remains to describe $P_m$. Each input node of $P_m$ (=an output node of $P_{m-1}$) has out-degree = 1. In any case the outgoing edge leads to a node $a$ if the last (the rightmost) digit of the transition after $C_1, ..., C_{m-1}$ is zero or to a node $b$ in the other case. Starting with $a, b$ $P_m$ consists of levels $L_i$, $i = 1, ..., l_m$. (We know that $l_m = m$.) Each level $L_i$ has two input nodes and two output nodes. Each input node of $L_i$ represents one value of the last digit of the sum of the transition after $C_1, ..., C_{m-1}$ and of $c_{1m}, ..., c_{i-1m}$. The output nodes of $L_i$ are stucked with the input nodes of $L_{i+1}$. The output nodes of $L_m$ are the sinks of $P$.

Now we see that $P$ will compute $f$.

At each input node of $L_i$ it is necessary to find out the value of $c_{im}$. This is ensured by two full trees of the depth 2 with roots in the input nodes. In both trees at the root the variable $x^{m-i+1}$ is tested and at the both successors of the root we test the variable $y^i$. With exception of $L_1$ the leaves of both trees are stucked into the output nodes of $L_i$.

Inside of $L_1$ firstly, in the tree with the root in $a$ ($b$, resp.) we stick the branches $x^0 = 0, y^1 = 0$ and $x^1 = 1, y^1 = 0$ to a node $c$ ($d$, resp.). Then we give edges leading from $c, d$ and the remaining leaves of both trees to the output nodes of $L_1$.

We see hat $|P| \leq O(n^3)$. It remains to prove that $P$ is a gentle branching program.

Let us perform the distribution on $P$. For all inputs no bit is double-crossed in $P_1, ..., P_{m-1}$ since each bit is tested in $P_m$. The nodes $c, d$ in $L_1$ of $P_m$ are
the first candidates for nodes with a proper class of the distribution.

**Proposition 7.5** Each input \( xy \) going through \( c \) (\( d \), resp.) has a double-cross (\( \# \)) on \( x^m \) and only on \( x^m \) at \( c \) (\( d \), resp.).

**Proof:** W.l.o.g. let \( xy \) be an input going through \( c \). At \( a \) there is the first test on \( x^m \). Hence also the input \( \bar{x}y \) goes through \( c \) where \( x \) and \( \bar{x} \) differ only on \( x^m \) (\( y^1 \) is non-crossed, \( y^1 = 0 \)). Since the test at \( a \) is the last test on \( x^m \), after \( c \) \( xy \) and \( \bar{x}y \) never branch and therefore they have a \( \# \) on \( x^m \). After \( c \) each input goes through tests on all remaining bits different from \( x^m \), \( y^1 \). Hence on these bits there is no \( \# \) at \( c \).

\( \square \)

\( 2^{n-1} \) inputs go through nodes \( c, d \). Hence at least one of the classes of the distribution at \( c \) and at \( d \) is a significant one. W.l.o.g. we assume that there is a significant class \( F \) at \( c \). We see that \( F \) is also an \( fd \)-subclass and therefore it suffices to prove that \( F \) satisfies the conditions (R1),(R2),(R3) from the definition of gentle branching programs.

Ad (R1). Let us consider the tree \( S_F \) induced by \( F \) at \( c \).

**Proposition 7.6** Let \( xy \) be an input from \( F \) following a branch \( b \) of \( S_F \). Then all bits different from \( x^1, x^2, x^3, x^m, y^1, y^{m-2}, y^{m-1}, y^m \) are crossed in \( w(xy, c, F, P) \).

**Proof:** By contradiction. Let \( e \) be a non-crossed bit of \( w(xy, c, F, P) \) different from \( x^1, ..., y^m \). We shall construct another input \( x'y' \) such that a) \( x'y' \in F \), b) in \( S_F \) \( x'y' \) follows \( xy \) till the test on \( e \) where \( xy, x'y' \) will branch. (Hence \( xy \) will have a cross on \( e \).) This will be our contradiction.

We choose \( x^m \) arbitrarily, we put \( y^1 = d_f \) 0. Then (in \( P \)) we follow \( xy \) from \( c \) to the test on \( e \) - this path defines the value of pairs \( x'^{m-1}, y'^2; x'^{m-2}, y'^3 \); \( ... \). In \( x'y' \) on \( e \) we give the opposite value (as in \( xy \)). In any case \( x'^1, x'^2, x'^3, y'^m, y'^{m-1}, y'^{m-2} \) remain undefined. If some other bits of \( x'y' \) are also undefined we define them arbitrarily. Further we define \( x'^1 = 1, x'^2 = 1, x'^3, y'^m \) arbitrarily. \( y'^{m-1} \) and \( y'^{m-2} \) are defined in such a way that the last digit of the transition after \( C_1, ..., C_{m-1} \) (on \( x'y' \)) is equal to 0.

We can argue that this is possible as follows: The starting point is to choose \( y'^{m-1} = y'^{m-2} = 0 \). If the last digit of the transition for \( C_m \) is equal to one we change \( y'^{m-1}, y'^{m-2} \). We know that \( y'^{m-1} \) is taken into account only in \( C_{m-1} \) (the corresponding bit is \( x'^1 \)) and that \( y'^{m-2} \) is taken into account only in \( C_{m-2} \) (the corresponding bit is \( x'^2 \)) in \( C_{m-2} \) (the corresponding bit is \( x'^1 \)). Let \( z' = d_f x' \times y' \). If under the condition \( y'^{m-1} = y'^{m-2} = 0 \) we have \( z'^{m-1} = 1 \) we change \( y'^{m-1} \) to 1. If \( z'^{m-1} = 0 \) and \( z'^{m-2} = 0 \) we put \( y'^{m-1} = y'^{m-2} = 1 \). If \( z'^{m-1} = 0 \) and \( z'^{m-2} = 1 \) we put \( y'^{m-1} = 0, y'^{m-2} = 1 \).
Hence $x'y$ goes through $a$ and moreover since $y'^1 = 0$ $x'y'$ goes through $c$ and it has $\neq$ on $x^m$. Therefore $x'y \in F$ - the condition a) from our plan of the proof is satisfied. Also b) is satisfied. A contradiction.

From the proposition it follows that the only candidates for non-crossed bits are $x^1, x^2, x^3, y^{m-2}, y^{m-1}, y^m$. Hence (R1) is satisfied.

Ad (R2). Let us choose $h \in H_F, t \in T_F$. Let us consider $\text{comp}(ht)$ in $P$. We shall prove that $[ht]$ goes through $c$ and after $c$ the computation depends only on $t$. Since $t \in T_F$ there is an input $xy \in F$ such that $xy \upharpoonright t_F = t$ (where $t_F$ contains all bits except of $x^m$). The input $xy$ goes from the source to $a$. Before $a$ there is no test on $x^m$. Hence $[ht]$ follows the same path from the source to $a$. Since $xy \in F$ $y^1 = 0$ holds - therefore also $[ht]$ goes to $c$. After $c$ only the bits from $t_F$ are tested, hence the computation depends only on $t$. We see that (R2) is satisfied.

Ad (R3) The bit which is double-crossed at $c$ is tested before $c$ at $a$ for all inputs from $F$. Hence (R1) is satisfied.

We see that $P$ is a gentle branching program. The Theorem is proven.

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