# Multihead one-way finite automata* 

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Communicated by A. Meyer
Received May 1988
Revised April 1989

## Abstract

Kutylowski, M., Multihead one-way finite automata, Theoretical Computer Science 85 (1991) 135-153.

We consider one-way non-sensing multihead finite automata. Let

$$
P_{m}=\left\{1^{a_{1}} * 1^{a_{2}} * \cdots * 1^{a_{m}} \# 1^{a_{m}} * 1^{a_{m-1}} * \cdots * 1^{a_{1}}: a_{1}, \ldots, a_{m} \in \mathbb{N}\right\} .
$$

We show that no $k$-head automaton can recognize the language $P_{m}$ if $m>\frac{1}{2} k^{3}$. It partially confirms the conjecture of Rosenberg. It shows that the languages $P_{m}$ and the languages $L_{m}$, where

$$
L_{m}=\left\{w_{1} * w_{2} * \cdots * w_{m} \# w_{m} * w_{m-1} * \cdots * w_{1}: w_{1}, \ldots, w_{m} \in\{0,1\}^{*}\right\}
$$

are of similar complexity for one-way multihead finite automata. We present a technique which can be used in some cases to estimate computational complexity of languages with respect to multihead automata.

## 1. Introduction

Multihead finite automata have been introduced in the early sixties by Piatkowski [11]. Obviously, multihead one-way finite automata (1-MFA) can recognize much more than regular languages and many authors have tried to characterize this class of languages. This problem received no satisfactory answer. For instance, the wellknown Pattern Matching Problem (see [4]) is still open. It concerns the language

$$
L_{s m}=\left\{w \# y w z: w, y, z \in\{0,1\}^{*}\right\} .
$$

[^0]It is unknown whether $I_{s m}$ can be recognized by some 1-MFA. There are some partial answers [9, 2], but it seems that we are still far from solving the problem. From the very beginning the problem of hierarchy with respect to the number of heads was considered for multihead automata. After many efforts [12-15, 3, 7] it was finally proved by Yao and Rivest [16] that " $k+1$ heads are better than $k$ heads". Their method was used afterwards by Hromkovič [5,6] to show that the class of languages recognized by $k$-head 1-MFA is not closed under many simple operations.

It turns out that the language $\left\{w \# w^{\mathrm{R}}: w \in\{0,1\}^{*}\right\}$ (where $w^{\mathrm{R}}$ stands for word $w$ written in the reverse order) cannot be recognized by any 1-MFA. However, various simpler versions of this language have been used to show several results about multihead automata. Consider e.g. the language

$$
L_{m}=\left\{w_{1} * w_{2} * \cdots * w_{m} \# w_{m} * w_{m-1} * \cdots * w_{1}: \forall j \leqslant m w_{j} \in\{0,1\}^{*}\right\} .
$$

Intuitively, this language is difficult for 1-MFA for the following reason. If automaton $M$ accepting language $L_{m}$ works on input $w_{1} * w_{2} * \cdots * w_{m} \# v_{m} * v_{m-1} * \cdots * v_{1}$, then it has to check that for each $j$ the blocks $w_{j}$ and $v_{j}$ are equal. A straightforward method to verify this is to put two heads at the beginnings of blocks $w_{j}$ and $v_{j}$ and then to move these heads simultaneously through words $w_{j}$ and $v_{j}$ checking if the corresponding symbols are equal. Note that each pair of heads can be used only once. Indeed, after checking $w_{j}$ and $v_{j}$ one head has only blocks $v_{j-1}, v_{j-2}, \ldots, v_{1}$ left to read, while the other one has only blocks $w_{j+1}, w_{j+2}, \ldots, w_{m}$ on its right side before the symbol \#. This old idea was used by Yao and Rivest [16] to show that $k$-head one-way automata can recognize languages $L_{m}$ only for $m \leqslant\binom{ k}{2}$. Their hierarchy theorem was just a simple corollary of this fact. To carry out the proof they had to use the fact that for a fixed $n$ there are about $2^{n}$ words of length $n$ in $L_{m}$, even if we fix the length of blocks $w_{j}$ and $v_{j}$. In this paper we try to answer if this argument must really be used. We consider relatively simple languages $P_{m}$, where

$$
P_{m}=\left\{1^{a_{1}} * 1^{a_{2}} * \cdots * 1^{a_{m}} \# 1^{a_{m}} * \cdots * 1^{a_{1}}: a_{1}, \ldots, a_{m} \in \mathbb{N}\right\} .
$$

Each block in $P_{m}$ is a block of 1's, so only its length must be checked against the length of the corresponding block. The proof used for languages $L_{m}$ does not work for languages $P_{m}$ for the reason that there are relatively few words in $P_{m}$ of a given length. The above languages were considered already by Rosenberg [13] in his attempt to prove the hierarchy theorem for 1-MFA.

In this paper we consider only non-sensing automata (a sensing multihead automaton can determine when coincidence of its heads occurs, non-sensing automata cannot detect whether their heads stay at the same place). We consider the following problem.

Problem 1.1. Given $k \in \mathbb{N}$. Find the maximal number $m$ such that language $P_{m}$ can be recognized by some $k$-head (non-sensing) 1-MFA.

We prove in Theorem 2.2 that for $m>k^{3} / 2$ there is no $k$-head 1 -MFA recognizing language $P_{m}$. On the other hand (see [16]), $k$-head automata are capable of recognizing languages $P_{m}$ for $m \leqslant\left(k^{2}-k\right) / 2$. We also prove (Theorem 3.1) that for each $k$ there
is a $k$-head automaton which accepts a language $P^{\prime} \subseteq P_{m}$, where $m \approx k^{3} / 24$. The language $P^{\prime}$ contains language $P_{m, c}$ which is a fairly large subset of $P_{m}$, namely,

$$
P_{m, c}=\left\{1^{a_{1}} * 1^{a_{2}} * \cdots * 1^{a_{m}} * 1^{a_{m}} * \cdots * 1^{a_{1}}: \forall k, l \leqslant m \quad a_{k} \leqslant c \cdot a_{l}\right\} .
$$

Anyway, the first-mentioned result shows that for non-sensing automata complexities of $P_{m}$ and $L_{m}$ are similar. The aim of this paper is not merely to give an answer to Problem 1.1, which is of rather technical importance. We have in mind a more general problem.

Problem 1.2. Let $L$ be a language of the form

$$
L=\left\{1^{f_{1}(n)} * 1^{f_{2}(n)} * \cdots * 1^{f_{m}(n)}: n \in \mathbb{N}\right\},
$$

where $f_{1}, \ldots, f_{m}$ are some functions over $\mathbb{N}$. Find a minimal number $k$ such that language $L$ can be recognized by a k-head (non-sensing) 1-MFA or show that $L$ cannot be recognized by such a device.

We do not know a complete answer to Problem 1.2. However, the methods used to solve Problem 1.1 can be easily adopted to provide a technical framework allowing to give answers to many subcases of Problem 1.2.

At this moment we have to mention the results obtained by Chrobak [1]. He considered very simple languages, namely,

$$
C_{n}=\left\{1^{x} 2^{i x}: i, x \in \mathbb{N}, 1 \leqslant i \leqslant n\right\},
$$

containing only two blocks of different symbols. He proved that for each $k$ there is an $n$ such that $C_{n}$ can be recognized by some $k$-head 1-MFA, but $C_{n+1}$ requires already $k+1$ heads. His proof was based on some geometric interpretation of the behaviour of 1-MFA. We follow this idea of such an interpretation. Also by applying some methods of this paper it is possible to determine how many heads are necessary to recognize language $C_{n}[8]$.

## 2. The upper bound

In this section we present our main result about recognizing languages $P_{m}$. From now on, by automata we mean multihead deterministic non-sensing one-way finite automata (1-MFA). For the sake of completeness we recall their definition taken from [7].

Definition 2.1. A one-way $k$-head deterministic non-sensing automaton is a device $M=\left\langle k, K, \Sigma, \delta, q_{0}, \$, F\right\rangle$, where $k \geqslant 1$ is the number of heads, $K, \Sigma$ and $F$ are finite sets of states, input symbols and accepting states $(F \subseteq K)$, respectively, $q_{0}$ is the initial state, $\$$ (not in $\Sigma$ ) is the right endmarker for the inputs, and $\delta$ is a mapping from $K \times(\Sigma \cup\{\$\})^{k}$ into $K \times\{0,1\}^{k}$. An input to $M$ is a string $a_{1} a_{2} \ldots a_{n}$ of symbols in $\Sigma$ delimited on the right end by the symbol $\$$. We can think of $a_{1} a_{2} \ldots a_{n} \$$ as written on
the tape (with each symbol occupying one tape square) and the heads moving left to right on the tape. One execution step of $M$ is described as follows. Let $M$ be in state $q$ with heads $H_{1}, \ldots, H_{k}$ scanning symbols $b_{1}, \ldots, b_{k}$ (in $\Sigma \cup\{\$\}$ ). Suppose $\delta\left(q, b_{1}, \ldots, b_{k}\right)=\left(p, d_{1}, \ldots, d_{k}\right)$. Then for each $j \leqslant k$, automaton $M$ moves head $H_{j}$ exactly $d_{j}$ squares to the right and enters state $p$ (we assume that no head scanning $\$$ can move further to the right). We say that a word $a_{1} a_{2} \ldots a_{n}$ is accepted or recognized by $M$ if when $M$ started on $a_{1} a_{2} \ldots a_{n} \$$ in state $q_{0}$ with all heads placed at $a_{1}$, after some number of steps it reaches an accepting state. Language $L$ is recognized by automaton $M$ if for every word $\boldsymbol{x}, \boldsymbol{x} \in L$ if and only if, $\boldsymbol{x}$ is accepted by $M$.

Theorem 2.2. Suppose $m>k^{3} / 2$. Then language $P_{m}$ cannot be recognized by any $k$-head (non-sensing) 1-MFA.

The rest of this section is devoted to a proof of this theorem. Assume that $M$ is a $k$-head automaton recognizing $P_{m}$. Consider input words of the form $\boldsymbol{x}=1^{a_{1}} *^{a_{2}} * \cdots * 1^{a_{2 m}}$. The subwords $1^{a_{1}}, 1^{a_{2}}, \ldots, 1^{a_{2 m}}$ are called blocks of $\boldsymbol{x}$, so $\boldsymbol{x}$ is built from $2 m$ such blocks. The $l$ th block $(l \leqslant 2 m)$ of $\boldsymbol{x}$ is denoted by $B_{l}$. Also the right endmarker $\$$ forms one additional block. Let the reading heads of $M$ be called $H_{1}$, $H_{2}, \ldots, H_{k}$.

Consider the computation on $\boldsymbol{x}$ performed by $\boldsymbol{M}$. It can be divided into several stages with each stage terminating when some head crosses a boundary between two adjacent blocks of $\boldsymbol{x}$. There are $k$ heads, so after at most $2 m k$ stages each computation must terminate. Without loss of generality we may assume that $M$ can reach a final state only when it enters a new stage of computation.

Now we shall construct $C_{M}$, a computation tree of $M$ with so-called distance functions and case conditions giving a full description of each possible computation performed by $M . C_{M}$ is a finite tree of height at most $2 m k$. For a node $a$ of $C_{M}$ let ht $(a)$ denote the height of $a$, i.e. $\mathrm{ht}(a)=l$ if $a$ is the $l$ th node on the path leading from the root of $C_{M}$ to node $a$. Each node of $C_{M}$ corresponds to the beginning of some stage of computation for some class of inputs. The root of $C_{M}$ corresponds to the beginning of the first stage (which is the same for all inputs). Its successors correspond to the second stages and so on: a node $a$ of $C_{M}$ corresponds to the stage ht $(a)$. With each node $a$ of $C_{M}$ we associate
(i) a state of $M$ denoted by $\operatorname{st}(a)$,
(ii) distance functions $f_{a, n, l}$ for $n \leqslant k, l \leqslant 2 m+1$,
(iii) case conditions.

Each case condition is an expression of the form

$$
h\left(x_{1}, x_{2}, \ldots, x_{2 m}\right)=0 \quad \text { or } \quad h\left(x_{1}, x_{2}, \ldots, x_{2 m}\right)>0,
$$

where the variables $x_{1}, x_{2}, \ldots, x_{2 m}$ stand for the lengths of blocks $B_{1}, B_{2}, \ldots, B_{2 m}$. We shall set $f_{a, n, l}\left(x_{1}, x_{2}, \ldots, x_{2 m}\right)$ to be the distance between head $H_{n}$ and the end of the $l$ th block of $1^{x_{1}} * \cdots * 1^{x_{2 m}}$ at the beginning of stage $\mathrm{ht}(a)$, if $H_{n}$ lies in this block. If $H_{n}$ does not lie in this block, then $f_{a, n, l}$ is undefined.

Let $p(a)$ denote the path connecting node $a$ with the root of $C_{M}$. Define $X_{a}$ to be the set of all words $\boldsymbol{x}, \boldsymbol{x}=1^{a_{1}} * 1^{a_{2}} * \cdots * 1^{a_{2 m}}$, such that for each case condition $g\left(x_{1}, \ldots, x_{2 m}\right)$ associated with a node lying on path $p(a), g\left(a_{1}, a_{2}, \ldots, a_{2 m}\right)$ holds. The key point of the construction of $C_{M}$ are the following properties.

Claim 2.3. (i) Suppose $a$ is a node of $C_{M}$. Then for every $\boldsymbol{x} \in X_{a}, \boldsymbol{x}=1^{a_{1}} *^{a_{2}} * \cdots * 1^{a_{2 m}}$, at the beginning of the stage corresponding to node a, automaton $M$ is in state st $(a)$ and for each $n \leqslant k$ the distance functions $f_{a, n, l}\left(a_{1}, a_{2}, \ldots, a_{2 m}\right)$ describe correctly the position of head $H_{n}$.
(ii) Each distance function is a "simple" expression (for a definition see below) and each case condition is of the form

$$
\xi\left(x_{1}, \ldots, x_{2 m}\right)=0 \quad \text { or } \quad \xi\left(x_{1}, \ldots, x_{2 m}\right)<c,
$$

where $\bar{\xi}$ is a simple expression.
Before the definition recall that $r(d, c)$ is the remainder on dividing $d$ by $c$.
Definition 2.4. (i) Remainder functions are defined inductively as follows. $R\left(x_{1}, \ldots, x_{2 m}\right)$ is a remainder function if

$$
R\left(x_{1}, \ldots, x_{2 m}\right)=c \cdot r\left(\sum_{j \leqslant n} \xi_{j}, e\right)
$$

where $c \in \mathbb{Z}, n, e \in \mathbb{N}$ and each $\xi_{j}$ is a remainder function, an integer or an integer multiplied by some $x_{i}, i \leqslant 2 m$.
(ii) We say that $\xi\left(x_{1}, \ldots, x_{2 m}\right)$ is a simple expression if it takes the form

$$
\xi\left(x_{1}, \ldots, x_{2 m}\right)=\left(\sum_{i \leqslant n} \xi_{i}\left(x_{1}, \ldots, x_{2 m}\right)\right) / d
$$

where $n, d \in \mathbb{N}$ and for each $i \leqslant n$ expression $\xi_{i}$ takes one of the possible forms: $c \cdot x_{j}, c$ or $R\left(x_{1}, \ldots, x_{2 m}\right)$, where $c \in \mathbb{Z}, j \leqslant 2 m$ and $R$ is a remainder function.

To define $C_{M}$ we construct inductively a finite sequence of trees $T_{0} \subseteq T_{1} \subseteq T_{2} \subseteq \ldots$ The last tree in this sequence is $C_{M}$. Each $T_{i}$ is an initial subtree of $C_{M}$, i.e. if $a \in C_{M}$ and $a \in T_{i}$, then $p(a) \subseteq T_{i}$. Therefore, to prove Claim 2.3 for $C_{M}$ it suffices to show it for each $T_{i}$, i.e. with $C_{M}$ replaced by $T_{i}$.
$T_{0}$ is a tree consisting only of $r$, the root of $C_{M}$. Clearly, st $(r)$ is the initial state of $M$, the set of the case conditions associated with $r$ is empty, the distance functions are defined as follows.

$$
\begin{array}{ll}
f_{r, n, 1}\left(x_{1}, \ldots, x_{2 m}\right)=x_{1} & \text { for } n \leqslant k \\
f_{r, n, j}\left(x_{1}, \ldots, x_{2 m}\right) & \text { is undefined for } j>1
\end{array}
$$

Now we assume that we have constructed $T_{i}(i \geqslant 1)$. If for every leaf $a$ of $T_{i}$ the state st $(a)$ is final, then the construction of $C_{M}$ is finished and $C_{M}=T_{i}$. Otherwise, we single out a leaf $a$ such that st $(a)$ is not a final state. Tree $T_{i+1}$ will be constructed by adding some number of successors to node $a$. Throughout the construction we shall consider exclusively input words from $X_{a}$. Let $n=h t(a)$. If $\boldsymbol{x} \in X_{a}, \boldsymbol{x}=1^{a_{1}} * 1^{a_{2}} * \cdots * 1^{a_{2 m}}$, then $M$ enters stage $n$ of computation over $\boldsymbol{x}$ in state $\operatorname{st}(a)$ with the head positions determined by values $f_{a, t, j}\left(a_{1}, a_{2}, \ldots, a_{2 m}\right)$ for $t \leqslant k, j \leqslant 2 m$. We do not know which head first reaches a boundary between blocks at the end of stage $n$. There are many possibilities, each of them giving rise to a different group of successors of $a$. So assume that $S$ is a set of heads and for each $j$, head $H_{j}$ crosses a boundary between blocks at the end of stage $n$ iff $H_{j} \in S$. We shall describe all successors of node $a$ corresponding to this situation. Let $H_{t} \in S$; head $H_{t}$ crosses the boundary between blocks $B_{l}$ and $B_{l+1}$ at the end of stage $n$. Hence, during this stage $H_{t}$ makes $f_{a, i, l}=f_{a, t, l}\left(a_{1}, \ldots, a_{2 m}\right)$ moves. Note that since each block contains only 1's (except the first symbol and the block $\$$ ), after some initial $p_{0}$ steps the behaviour of $M$ in stage $n$ becomes periodic. Of course, $f_{a, t, l}$ can be too small to reach the first period cycle. For each $p \leqslant p_{0}$ there is a successor of $a$ in $T_{i+1}$ describing the situation when stage $n$ terminates after exactly $p$ machine steps. So assume that machine $M$ makes $p$ steps ( $p \leqslant p_{0}$ ) and stage $n$ terminates. We describe the corresponding node $a^{\prime}$ of $T_{i+1}$ added to $T_{i}$. For each head $H_{s}$ of $M$ let $e_{s}$ be the number of moves made by $H_{s}$ during those $p$ machine steps in stage $n$. Iet $e_{s}^{\prime}$ be the number of moves made by $H_{s}$ during $p-1$ such steps. The distance functions of $a^{\prime}$ are defined as follows.

For $s \notin S$

$$
\begin{array}{ll}
f_{a^{\prime}, s, j}=f_{a, s, j}-e_{s} & \text { if } f_{a, s, j} \text { is defined, } \\
f_{a^{\prime}, s, j} & \text { is undefined otherwise. }
\end{array}
$$

For $s \in S$ function $f_{a^{\prime}, s, j+1}$ is defined iff $f_{a, s, j}$ is defined and $f_{a^{\prime}, s, j+1}=x_{j+1}$. The case conditions associated with $a^{\prime}$ are defined as follows:

$$
\begin{array}{lll}
f_{a, s, j}-e_{s}>0 & \text { for } s \notin S, \\
f_{a, s, j}-e_{s}=0 \quad \text { and } \quad f_{a, s, j}-e_{s}^{\prime}>0 & \text { for } s \in S .
\end{array}
$$

In the above conditions $j$ is chosen such that $f_{a, s, j}$ is defined. Clearly, node $a^{\prime}$ satisfies Claim 2.3(i). For Claim 2.3(ii) note that a sum of a simple expression and a constant is a simple expression.

Now we have to consider the case when $f_{a, t, l}$ is big enough to reach the periodic behaviour during stage $n$. We have to determine the number of moves made by each head. Consider head $H_{s}$. Before $M$ reaches the first moment of the periodic part of the execution, $H_{t}$ makes some $c_{t}$ moves and $H_{s}$ makes some $c_{s}$ moves ( $c_{t}, c_{s}$ are constants). So there are still $f_{a, t, l}-c_{t}$ cells left for $H_{t}$ in block $B_{l}$. During each period cycle, $H_{t}$ makes $u_{t}$ moves while $H_{s}$ makes $u_{s}$ moves, for some constants $u_{t}, u_{s}$. So, to the end of the last fully executed cycle $H_{t}$ makes

$$
\left(f_{a, t, l}-c_{t}\right)-r\left(f_{a, t, l}-c_{t}, u_{t}\right)
$$

moves. It corresponds to $\left[\left(f_{a, t, t}-c_{t}\right) / u_{t}\right.$ ] full cycles (where $[p]$ stands for the integer part of $p$, the maximal $n \in \mathbb{Z}$ such that $n \leqslant p$ ). At the same time head $H_{s}$ makes

$$
\left[\left(f_{a, t, l}-c_{t}\right) / u_{t}\right] \cdot u_{s}
$$

moves. During the last, not fully executed cycle each head makes some constant number of moves depending only on $r\left(f_{a, t, l}-c_{t}, u_{t}\right)$. Machine $M$ needs some $p_{1}$ steps to execute the cycle. For each $p<p_{1}$ there is a node in $T_{i+1}$ corresponding to the situation when the last cycle terminates after $p$ initial steps. We fix $p<p_{1}$ and we describe the corresponding node $a^{\prime \prime}$ of $T_{i+1}$. Let $v_{s}$ be the number of moves made by $H_{s}$ during $p$ initial machine steps in the cycle, and let $v_{s}^{\prime}$ be the number of moves made during $p-1$ such steps. Finally, we can say that during stage $n$ head $H_{s}$ makes

$$
c_{s}+\left[\left(f_{a, t, l}-c_{t}\right) / u_{t}\right] \cdot u_{s}+v_{s}
$$

moves. New distance functions associated with $a^{\prime \prime}$ are defined as follows. For heads $H_{s} \in S$ it is simple since these heads enter new blocks. Let $s \notin S$. Then the new distance function $f_{a^{\prime \prime}, s, j}$ is undefined if $f_{a, s, j}$ is undefined and

$$
f_{a^{\prime \prime}, s, j}=f_{a, s, j}-c_{s}-\left[\left(f_{a, t, l}-c_{t}\right) / u_{t}\right] \cdot u_{s}-v_{s}
$$

otherwise. The case conditions associated with $a^{\prime \prime}$ are the following.
For each $s \notin S$ and $j$ such that $f_{a, s, j}$ is defined

$$
\begin{aligned}
& f_{a, s, j}-c_{\mathrm{s}}-\left[\left(f_{a, t, l}-c_{\mathrm{t}}\right) / u_{\mathrm{t}}\right] \cdot u_{\mathrm{s}}-v_{\mathrm{s}}=0, \\
& f_{a, s, j}-c_{\mathrm{s}}-\left[\left(f_{a, t, t, l}-c_{t}\right) / u_{t}\right] \cdot u_{s}-v_{s}^{\prime}>0
\end{aligned}
$$

For $s \in S$ and $j$ such that $f_{a^{\prime \prime}, s, j}$ is defined

$$
f_{a^{\prime \prime}, s, j}>0
$$

It follows from the construction that Claim 2.3(i) holds for $T_{i+1}$. So it remains to prove that each of the above-defined case condition and distance function is of the desired form. What we need is the following lemma.

Lemma 2.5. (i) A sum of simple expressions is a simple expression.
(ii) If $\xi$ is a simple expression and $c \in \mathbb{Z}$, then $c \cdot \xi$ is also a simple expression.
(iii) If $\xi$ is a simple expression, $c \in \mathbb{Z}, d \in \mathbb{N}$, then $[(\xi-c) / d]$ is a simple expression.

Proof. (i) and (ii) are obvious. For (iii) consider $\xi=\left(\sum \xi_{j}\right) / e$. Then

$$
\begin{align*}
{[(\xi-c) / d] } & =\frac{\xi-c-r(\xi-c, d)}{d} \\
& =\frac{\left(\sum \xi_{j}\right) / e-c-r\left(\left(\sum \xi_{j}\right) / e-c, d\right)}{d} \\
& =\frac{\sum \xi_{j}-c e-e \cdot r\left(\left(\sum \xi_{j}-c e\right) / e, d\right)}{d e} . \tag{2.1}
\end{align*}
$$

Now we have to prove that $e \cdot r\left(\left(\sum \xi_{j}-c e\right) / e, d\right)$ is of the appropriate form. For $\alpha, \beta, \sigma \in \mathbb{N}, \alpha / \beta \in \mathbb{N}$ we have $r(\alpha / \beta, \sigma)=r(\alpha, \beta \sigma) / \beta$. So

$$
e \cdot r\left(\left(\sum \xi_{j}-c e\right) / e, d\right)=e \cdot r\left(\sum \xi_{j}-c e, e d\right) / e=r\left(\sum \xi_{j}-c e, e d\right) .
$$

The last expression is a remainder function, so by $(2.1),[(\xi-c) / d]$ is a simple expression.

By Lemma 2.5, each constructed distance function and case condition has the desired form and Claim 2.3(ii) holds for the tree $T_{i+1}$.
We have just described the construction of tree $C_{M}$. Since each node of $C_{M}$ has finitely many successors and the height of $C_{M}$ is not greater than $2 m k, C_{M}$ is a finite tree.

While looking at the work of $M$ it is troublesome to consider all remainder functions which occur within the case conditions and the distance functions. We shall find a way to elude this difficulty. Let $\mathscr{R}$ be the set of all these remainder functions. Take $g \in \mathbb{N}$ such that if $R \in \mathscr{R}$, say $R\left(x_{1}, \ldots, x_{2 m}\right)=c \cdot r\left(\sum \xi_{j}, e\right)$, then $e$ divides $g$.

Definition 2.6. Let $X^{\prime}=\left\{1^{a_{1}} * 1^{a_{2}} * \cdots * 1^{a_{2 m}}: \forall i g \mid a_{i}\right\}$ and $P^{\prime}=P_{m} \cap X^{\prime}$.
Essentially, in the rest of this section we shall use only inputs from $X^{\prime}$. The reason for that will become clear when we formulate the following lemma.

Lemma 2.7. For each $R \in \mathscr{R}$ there is a constant $c_{R}$ such that for every $\boldsymbol{x} \in X^{\prime}$, $\boldsymbol{x}=1^{a_{1}} * 1^{a_{2}} * \cdots * 1^{a_{2 m}}$, we have $R\left(a_{1}, \ldots, a_{m}\right)=c_{R}$.

Proof. By induction on complexity of $R$ : Suppose $R\left(x_{1}, \ldots, x_{2 m}\right)=c \cdot r\left(\sum_{j}, e\right)$. If $\xi_{j}\left(x_{1}, \ldots, x_{2 m}\right)=d \cdot x_{j}$, then $\xi_{j}\left(a_{1}, \ldots, a_{2 m}\right)=d \cdot a_{j}$. But $e \mid g$ and $g \mid a_{j}$, so $e \mid d \cdot a_{j}$. If $\xi_{j}$ is a remainder function, then by the induction hypothesis $\xi_{j}\left(a_{1}, \ldots, a_{2 m}\right)$ has a constant value not depending on $a_{1}, \ldots, a_{2 m}$. Hence, each $\xi_{j}\left(a_{1}, \ldots, a_{2 m}\right)$ is either a constant not depending on $a_{1}, \ldots, a_{2 m}$ or a number divisible by $e$. So the lemma holds for $R$.

Definition 2.8. If $\phi$ is a distance function

$$
\begin{equation*}
\phi\left(x_{1}, \ldots, x_{2 m}\right)=\left(\sum_{i} \xi_{i}\left(x_{1}, \ldots, x_{2 m}\right)\right) / d \tag{2.2}
\end{equation*}
$$

then $\phi^{\prime}$, the reduced distance function, is obtained from $\phi$ by replacing in (2.2) each remainder function $\xi_{i}$ by the constant $c_{\xi_{i}}$ given by Lemma 2.7. Similarly, we define the reduced case conditions.

Lemma 2.9. If $\boldsymbol{x} \in X^{\prime}, \boldsymbol{x}=1^{a_{1}} * 1^{a_{2}} \ldots \ldots * 1^{a_{2 m}}$, then for each distance function $\phi$

$$
\phi\left(a_{1}, a_{2}, \ldots, a_{2 m}\right)=\phi^{\prime}\left(a_{1}, a_{2}, \ldots, a_{2 m}\right)
$$

and for each case condition $\psi$ we have

$$
\psi\left(a_{1}, a_{2}, \ldots, a_{2 m}\right) \Leftrightarrow \psi^{\prime}\left(a_{1}, a_{2}, \ldots, a_{2 m}\right) .
$$

Proof. It follows immediately from the definition.
The most important fact about the reduced distance functions and the reduced case conditions is that they involve only linear functions. It considerably simplifies the situation.

Assume that $l$ is a node of $C_{M}$. Then let $X_{I}^{\prime}=X_{l} \cap X^{\prime}$. Let $\Psi_{l}$ be the set of all reduced case conditions associated with nodes lying on path $p(l)$. Clearly,

$$
X_{l}^{\prime}=\left\{1^{a_{1}} * 1^{a_{2}} * \cdots * 1^{a_{2 m}} \in X^{\prime}: \forall \phi^{\prime} \in \Psi_{l} \quad \phi^{\prime}\left(a_{1}, a_{2}, \ldots, a_{2 m}\right)\right\} .
$$

Before we proceed we recall terminology of geometry which we shall use. If $\phi\left(x_{1}, \ldots, x_{n}\right)=\alpha_{1} x_{1}+\alpha_{2} x_{2}+\cdots+\alpha_{n} x_{n}$ (where $\alpha_{1}, \ldots, \alpha_{n} \in \mathbb{Q}$ ) and $c \in \mathbb{Q}$, then the set $\left\{\boldsymbol{x} \in \mathbb{Q}^{n}: \phi(\boldsymbol{x})=c\right\}$ is called a hyperplane in $\mathbb{Q}^{n}$. The sets $\left\{\boldsymbol{x} \in \mathbb{Q}^{n}: \phi(\boldsymbol{x})>c\right\}$ and $\left\{\boldsymbol{x} \in \mathbb{Q}^{n}\right.$ : $\phi(x) \geqslant c\}$ are called halfspaces with the edge $\left(x \in \mathbb{Q}^{n}: \phi(x)=c\right\}$. An intersection of a finite number of halfspaces is a polyhedron. Hence, for each polyhedron $U$ there are $\phi_{1}, \phi_{2}, \ldots, \phi_{j}$, linear combinations of $x_{1}, \ldots, x_{n}$, and $c_{1}, \ldots, c_{n} \in \mathbb{Q}$ such that $U=\left\{\boldsymbol{x} \in \mathbb{Q}^{n}\right.$ : $\left.\forall i \leqslant j \phi_{i}(\boldsymbol{x}) \$_{i} c_{i}\right\}$, where each $\$_{i}$ is $<$ or $\leqslant U$ is called a layer if there is $\phi\left(x_{1}, \ldots, x_{n}\right)$, a linear combination of $x_{1}, \ldots, x_{n}$, and $c_{1}, c_{2} \in \mathbb{Q}$ such that $U=\left\{\boldsymbol{x} \in \mathbb{Q}^{n}: c_{1} \leqslant \phi(\boldsymbol{x}) \leqslant c_{2}\right)$. So a layer is simply a set of points which lie between two parallel hyperplanes. A polygon is a polyhedron of dimension 2. A face of a polyhedron $U$ is either the empty set or a polyhedron obtained by replacing some of the inequalities that define $U$ with equations. A proper face of $U$ is a face not equal to the empty set or $U$. A maximal proper face of $U$ is called a facet of $U$. For $S \subseteq \mathbb{Q}^{n}$ the affine hull of $S$, aff $(S)$, is the set of all $z \in \mathbb{Q}^{n}$ which can be expressed as $z=\sum_{x \subset S^{\prime}} \lambda_{x} \cdot x$ satisfying $\sum_{x \subset S^{\prime}} \lambda_{x}=1$ for some finite $S^{\prime} \subseteq S$. $S$ is an affine subspace of $\mathbb{Q}^{n}$ if $\operatorname{aff}(S)=S$. We define the dimension $\operatorname{dim} S$ of $S \subseteq \mathbb{Q}^{n}$ to be the dimension of aff $(S)$, its affine hull.

Each word $1^{a_{1}} * 1^{a_{2}} * \cdots * 1^{a_{2 m}}$, corresponds to the string $\left(a_{1}, a_{2}, \ldots, a_{2 m}\right) \in \mathbb{Q}^{2 m}$. For several reasons it will be more convenient to consider elements of $\mathbb{Q}^{2 m}$ rather than words. This enables us to use simple geometric techniques. For that reason we shall identify word $1^{a_{1}} * \cdots * 1^{a_{2 m}}$ with the point $\left(a_{1}, \ldots, a_{2 m}\right) \in \mathbb{Q}^{2 m}$. Also the sets $X^{\prime}, X_{l}^{\prime}, P^{\prime}$ shall be treated as subsets of $\mathbb{Q}^{2 m}$.

Let $V=\left\{\boldsymbol{x} \in \mathbb{Q}^{2 m}: \forall i x_{i}=x_{2 m-i+1}\right\}$. Clearly, $P^{\prime}=V \cap X^{\prime}$ and $V$ is an affine subspace of $\mathbb{Q}^{2 m}$ of dimension $m$. Recall that $X_{l}^{\prime}=\left\{\boldsymbol{x} \in X^{\prime}: \forall \phi \in \Psi_{l} \phi(\boldsymbol{x})\right\}$. Define

$$
U_{I}=\left\{\boldsymbol{x} \in \mathbb{Q}^{2 m}: \forall \phi \in \Psi, \phi(\boldsymbol{x})\right\} .
$$

Then, obviously, $X_{l}^{\prime}=U_{l} \cap X^{\prime}$. Recall that each $\phi \in \Psi_{l}$ defines in $\mathbb{Q}^{2 m}$ a hyperplane or a halfspace depending on whether $\phi$ is an equation or an inequality. Hence, $U_{l}$ is a polyhedron in $\mathbb{Q}^{2 m}$. Automaton $M$ recognizes language $P_{m}$, so

$$
\begin{aligned}
& X^{\prime}=\bigcup\left\{X_{l}^{\prime}: l \text { is a leaf of } C_{M}\right\}, \\
& P^{\prime}=\bigcup\left\{X_{l}^{\prime}: l \text { is an accepting leaf of } C_{M}\right\} .
\end{aligned}
$$

Polyhedrons $U_{l}$ might not cover $\mathbb{Q}^{2 m}$ but by the first of the above equalities they cover $X^{\prime}$. Now we shall seek $l$, an accepting leaf of $C_{M}$, such that $U_{l}$ is large enough for our purposes. Before that we must prove some auxiliary facts of geometry.

Lemma 2.10. Suppose $g \in \mathbb{N}$. Take a finite set of layers in $\mathbb{Q}^{n}$, say $F_{1}, F_{2}, \ldots, F_{t}$. Then there is a point $\boldsymbol{x} \in \mathbb{Q}^{n} \backslash \bigcup_{i} F_{i}$ such that $\forall i \leqslant n g \mid x_{i}$.

Proof. By simple induction on $n$ : For $n=1$ the lemma is obviously true, so assume that $n>1$. Consider hyperplanes $W_{i}: W_{i}=\left\{x \in \mathbb{Q}^{n}: x_{1}=g \cdot i\right\}$. If the edges of $F_{j}$ are not parallel to hyperplane $\left\{\boldsymbol{x} \in \mathbb{Q}^{n}: x_{1}=0\right\}$, then $F_{j} \cap W_{i}$ is a layer in $W_{i}$ for each $i$. If the edges of $F_{j}$ are parallel to this hyperplane, then for almost all $i$ sets $F_{j}$ and $W_{i}$ are disjoint. So there is $i \in \mathbb{N}$ such that the sets $W_{i} \cap F_{1}, W_{i} \cap F_{2}, \ldots, W_{i} \cap F_{t}$ are layers in $W_{i}$. Hyperplane $W_{i}$ is isomorphic to $\mathbb{Q}^{n-1}$, so by the induction hypothesis there is $\boldsymbol{x} \in W_{i} \backslash \bigcup\left(W_{i} \cap F_{i}\right)=W_{i} \backslash \bigcup F_{i}$ such that $g \mid x_{i}$ for $i=2,3, \ldots, n$. But $\boldsymbol{x} \in W_{i}$ so $g \mid x_{1}$.

Let $\mathbb{Q}_{+}^{n}=\left\{\boldsymbol{x} \in \mathbb{Q}^{n}: \forall i \leqslant n x_{i} \geqslant 0\right\}$. Since $X^{\prime} \subseteq \mathbb{Q}_{+}^{2 m}$, we shall virtually stay in $\mathbb{Q}_{+}^{2 m}$.
Lemma 2.11. Lemma 2.10 holds also if we replace $\mathbb{Q}^{n}$ by $\mathbb{Q}^{n}+$.
Proof. $\mathbb{Q}^{n}$ is a union of finitely many subsets isomorphic to $\mathbb{Q}^{n}+$. So if we could cover $\mathbb{Q}_{+}^{n}$ by finitely many layers we could do the same with $\mathbb{Q}^{n}$.

Definition 2.12. Let $L$ be a polyhedron, say $L=\bigcap_{i \in J}\left\{\boldsymbol{x}: \phi_{i}(x) \$_{i} c_{i}\right\}$ where each $\$_{i}$ is $\leqslant$ or $<$. If $z \in L$ then by $\mathscr{L}(L, z)$ we mean the polyhedron (Fig. 1)

$$
\mathscr{L}(L, z)=\bigcap_{i \in J}\left\{x: \phi_{i}(x) \geqslant \phi_{i}(z)\right\} .
$$



Fig. 1.

Lemma 2.13. If $L$ is a polyhedron, $z \in L$, then
(i) $\mathscr{L}(L, z) \subseteq L$.
(ii) If $\mathscr{L}(L, z)$ is not the point $z$ itself, then

$$
\begin{equation*}
\mathscr{L}(L, z)=\bigcup\{p: p \text { is a halfine with the end } z, p \subseteq L\} . \tag{2.3}
\end{equation*}
$$

(iii) If $z, s \in L$ then $\mathscr{L}(L, s)$ is equal to $\mathscr{L}(L, z)$ translated by vector $\overrightarrow{\boldsymbol{z}}$.
(iv) If $\mathscr{L}(L, z) \subseteq V$ and $\operatorname{dim} \mathscr{L}(L, z)=m$, then $\mathscr{L}(L, z) \cap X^{\prime} \neq \emptyset$.

Proof. Let $L$ be defined as in Definition 2.12. Part (i) is obvious. For part (ii) we show first the inclusion $\subseteq$. Suppose $y \in \mathscr{L}(L, z), y \neq z$. Take the halfline $p$ with the end $z$ such that $y \in p$. Consider function $\phi_{i}$ on $p$. It is linear. We have $\phi_{i}(y) \geqslant \phi_{i}(z)$, so $\phi_{i}$ is not decreasing along $p$. So $\phi_{i}(\boldsymbol{u}) \geqslant \phi_{i}(\boldsymbol{z})$ for every $\boldsymbol{u} \in p$. It holds for every $i$ so $p \subseteq \mathscr{L}(L, z) \subseteq L$. Hence, $p$ witnesses that $y$ is an element of the union on the right side of (2.3). For the inclusion $\supseteq$ assume that $p$ is a halfline with the end $z, p \subseteq L$. Consider function $\phi_{i}$ on $p$. It is linear and has values not smaller than $c_{i}$. So function $\phi_{i}$ cannot decrease along $p$. IIence, its values are not smaller than $\phi_{i}(z)$. So we get $p \subseteq \mathscr{L}(L, z)$.

For (iii) we show first that $\mathscr{L}(L, z)+\vec{z} \subseteq \mathscr{L}(L, s)$. Take $\boldsymbol{u} \in \mathscr{L}(L, z)$. Then $\phi_{i}(\boldsymbol{u}) \geqslant \phi_{i}(z)$ for each $i$. Note that $\phi_{i}(\boldsymbol{u}+\overrightarrow{\boldsymbol{z} s})-\phi_{i}(\boldsymbol{u})=\phi_{i}(\boldsymbol{s}) \quad \phi_{i}(z)$, because $\phi_{i}$ is a linear function. So

$$
\begin{aligned}
& \left(\phi_{i}(\boldsymbol{u}+\overrightarrow{z s})-\phi_{i}(\boldsymbol{u})\right)+\phi_{i}(\boldsymbol{u}) \geqslant\left(\phi_{i}(s)-\phi_{i}(z)\right)+\phi_{i}(z), \\
& \phi_{i}(\boldsymbol{u}+\overrightarrow{\boldsymbol{z s}}) \geqslant \phi_{i}(\boldsymbol{s}) .
\end{aligned}
$$

So $(\boldsymbol{u}+\overrightarrow{\boldsymbol{z}}) \in \mathscr{L}(L, \boldsymbol{s})$. Hence, $\mathscr{L}(L, z)+\overrightarrow{z s} \subseteq \mathscr{L}(L, s)$. In the same way we get $\mathscr{L}(L, s)+\vec{s} \subseteq \mathscr{L}(L, z)$. So $\mathscr{L}(L, s)+\vec{s} \vec{z}+\vec{z} \subseteq \mathscr{L}(L, z)+\vec{z}$. Then $\mathscr{L}(L, s) \subseteq \mathscr{L}(L, z)+\vec{z}$.

For (iv) note that, since $\operatorname{dim}(\mathscr{L}(L, z))=m$ and $\mathscr{L}(L, z)$ is a polyhedron in $V$, there is a point $s$ which is an interior point of $\mathscr{L}(L, z)$ in the sense of topology of $V$. Then $\phi_{i}(s)>\phi_{i}(z)$ for each $i$ (otherwise, $s$ would lie on one of the facets of $\mathscr{L}(L, z)$ ). By (ii), there is a halfline $p$ beginning with $z$ and containing point $s$. Function $\phi_{i}$ grows to infinity on $p$. Since $\phi_{i}$ 's are linear functions, we can find a number $\eta \in \mathbb{Q}, \eta>0$, such that for each pair of points $\boldsymbol{x}, \boldsymbol{y},\left|\phi_{i}(\boldsymbol{x})-\phi_{i}(\boldsymbol{y})\right| \leqslant \eta \cdot d(\boldsymbol{x}, \boldsymbol{y})$, where $d(\boldsymbol{x}, \boldsymbol{y})$ denotes the distance between $\boldsymbol{x}$ and $\boldsymbol{y}$. Thercforc, we can find a point $\boldsymbol{u} \in p$ such that a ball $K$ in $V$ with center at $\boldsymbol{u}$ and radius $\frac{1}{2} g \sqrt{m}$ is a subset of $\mathscr{L}(L, z)$. Indeed, it suffices to take $\boldsymbol{u}$ such that $\phi_{i}(\boldsymbol{u})>\phi_{i}(z)+\frac{1}{2} \eta \cdot g \sqrt{m}$ for each $i$. If $\boldsymbol{x} \in K$, then

$$
\begin{aligned}
\phi_{i}(\boldsymbol{x}) & =\phi_{i}(\boldsymbol{u})+\left(\phi_{i}(\boldsymbol{x})-\phi_{i}(\boldsymbol{u})\right) \geqslant \phi_{i}(\boldsymbol{u})-\eta \cdot d(\boldsymbol{x}, \boldsymbol{u}) \\
& \geqslant \phi_{i}(\boldsymbol{z})+\frac{1}{2} \eta \cdot g \sqrt{m}-\frac{1}{2} \eta \cdot g \sqrt{m}=\phi_{i}(z) .
\end{aligned}
$$

So $x \in \mathscr{L}(L, z)$. Each ball in $V$ of radius $\frac{1}{2} g \sqrt{m}$ contains a point from $X^{\prime}$. So the lemma follows.

Now consider $l_{0}, l_{1}, \ldots, l_{t}$, all accepting leaves of $C_{M}$. We may assume that each $\Psi_{l_{i}}$, the set of reduced case conditions leading to $l_{i}$, contains the inequalities $x_{j} \geqslant 0$ for $j \leqslant 2 m$. Consider now polyhedrons $W_{i}$, where $W_{i}=V \cap U_{l_{i}}$. Of course, $V \cap \mathbb{Q}_{+}^{2 m}$ is
isomorphic to $\mathbb{Q}_{+}^{m}$ and $V \cap \mathbb{Q}_{+}^{2 m} \supseteq \bigcup_{i} W_{i} \supseteq V \cap \mathbb{Q}_{+}^{2 m} \cap X^{\prime}$. For each $W_{i}$ fix a point $\boldsymbol{w}_{i} \in W_{i}$.

Claim 2.14. For some i polyhedron $\mathscr{L}\left(W_{i}, w_{i}\right)$ has dimension $m$.
Proof. Assume the converse. So each $\mathscr{L}\left(W_{i}, \boldsymbol{w}_{i}\right)$ is a subset of some layer. On the other hand, by the definition of $\mathscr{L}\left(W_{i}, \boldsymbol{w}_{i}\right)$, the set $W_{i} \backslash \mathscr{L}\left(W_{i}, \boldsymbol{w}_{i}\right)$ is included in a finite union of layers. Consequently, we can cover $V \cap \mathbb{Q}_{+}^{2 m} \cap X^{\prime}$ by a finite union of layers in $V \cap \mathbb{Q}_{+}^{2 m}$. This contradicts Lemma 2.11.

By the above claim we may assume that polyhedron $\mathscr{L}\left(W_{0}, \boldsymbol{w}_{0}\right)$ has dimension $m$. By Lemma 2.13 there is $\boldsymbol{x}_{0} \in W_{0} \cap X^{\prime}$. Then again by Lemma 2.13, $\mathscr{L}\left(W_{0}, \boldsymbol{x}_{0}\right)$ has dimension $m$. Leaf $l_{0}$ is the leaf $l$ we have been looking for, for which $U_{l}$ is "large enough".

Polyhedron $U_{l_{0}}$ is defined by the reduced case conditions from $\Psi_{l_{0}}$. We split $\Psi_{l_{0}}$ into two subsets $\Phi_{1}$ and $\Phi_{2}$. For $\phi \in \Psi_{l_{0}}$ we put $\phi$ into $\Phi_{1}$ if either the hyperplane defined by $\phi$ (if $\phi$ is an equation) is parallel to $V$ or the halfspace defined by $\phi$ has the edge parallel to $V$. Otherwise, $\phi$ is in $\Phi_{2}$. Let

$$
V_{i}=\left\{\boldsymbol{x} \in \mathbb{Q}^{2 m}: \forall \phi \in \Phi_{i} \phi(\boldsymbol{x})\right\}
$$

for $i=1$, 2. Obviously, $U_{l_{0}}=V_{1} \cap V_{2}$ and
(i) all facets of $V_{1}$ are parallel to $V$,
(ii) no facet of $V_{2}$ is parallel to $V$.

Before we proceed let us notice the following property.
Property 2.15. If a halfine $p \subseteq \mathbb{Q}^{2 m}$ contains a point from $X^{\prime}$, then it contains infinitely many of them, each two subsequent points staying at a constant distance (depending on $p$ ).

The property easily follows from the fact that we are working in $\mathbb{Q}^{2 m}$, not in $\mathbb{R}^{2 m}$.
Lemma 2.16. $V_{1} \cap X^{\prime} \subseteq V$.
Proof. We know that $V_{1} \cap V_{2} \cap X^{\prime}=U_{l_{0}} \cap X^{\prime} \subseteq V$. Assume that there is a point $\boldsymbol{x} \in\left(V_{1} \cap X^{\prime}\right) \backslash V$.

Claim 2.17. There is a halfine $p \subseteq \mathscr{L}\left(W_{l_{0}}, \boldsymbol{x}_{0}\right)$ with the beginning $\boldsymbol{x}_{0}$ not parallel to any facet of $V_{2}$.

Proof of Claim 2.17. To show the claim note that for each facet $F$ of $V_{2}$ the set $\left\{\boldsymbol{y} \in \mathscr{L}\left(W_{l_{0}}, \boldsymbol{x}_{0}\right): \overrightarrow{\boldsymbol{x}_{0} \boldsymbol{y}}\right.$ is parallel to $\left.F\right\}$ has dimension less than $m$. A finite union of sets of dimension less than $m$ cannot cover polyhedron $\mathscr{L}\left(W_{l_{0}}, x_{0}\right)$ of dimension $m$. Hence, there is a point $\boldsymbol{y} \in \mathscr{L}\left(W_{l_{0}}, \boldsymbol{x}_{0}\right)$ such that vector $\overrightarrow{\boldsymbol{x}_{0} \boldsymbol{y}}$ is not parallel to any facet of $V_{2}$. Take $p$ to be the halfline beginning with $\boldsymbol{x}_{0}$ and containing point $\boldsymbol{y}$.


Fig. 2. The situation on plane $q$.

Proof of Lemma 2.16 (conclusion): Let $q$ be the plane containing halfline $p$ and point $\boldsymbol{x}$ (Fig. 2). Then, by Claim $2.17 V_{2} \cap q$ is a polygon with edges not parallel to $p$. Since $p \subseteq U_{l_{0}}, p \subseteq U_{l_{0}} \cap q \subseteq V_{2} \cap q$. Let $r$ be the line parallel to $p$ containing point $\boldsymbol{x}$. It is geometrically evident that $\left(V_{2} \cap q\right) \cap r$ is a halfline, say $s$. Since $x \in r$, it follows from Property 2.15 that $s$ contains infinitely many points from $X^{\prime}$. On the other hand, line $r$ is parallel to $p$ so $r$ is parallel to $V$. Polyhedron $V_{1}$ has facets parallel to $V$, line $r$ contains point $\boldsymbol{x}$ from $V_{1}$, so $r \subseteq V_{1}$. Hence, $s \subseteq V_{1} \cap V_{2}$. We have noticed previously that $\boldsymbol{s} \cap X^{\prime} \neq \emptyset$. So there is a point $\boldsymbol{x}^{\prime} \in \boldsymbol{s} \cap X^{\prime} \subseteq V_{1} \cap V_{2} \cap X^{\prime}=X_{l_{0}}^{\prime}$. Lcaf $l_{0}$ is accepting, so $\boldsymbol{x}^{\prime} \in V$. We know that $x^{\prime} \in V, p \subseteq V$ and $V$ is an affine space. So $q \subseteq V$ and hence, $\boldsymbol{x} \in V$ contrary to the assumption about $\boldsymbol{x}$.

Polyhedron $V_{1}$ has interesting properties: $V \subseteq V_{1}$ and $V_{1} \cap X^{\prime} \subseteq V$. Let $g_{1}, g_{2}, \ldots, g_{n}$ be all reduced case conditions defining $V_{1}$ (i.e. $\Phi_{1}=\left\{g_{1}, g_{2}, \ldots, g_{\eta}\right\}$ ).

Lemma 2.18. Each $g_{j}$ is a condition of the form

$$
h_{j}\left(x_{1}, \ldots, x_{2 m}\right)<c \quad \text { or } \quad h_{j}\left(x_{1}, \ldots, x_{2 m}\right)=c^{\prime},
$$

where $c^{\prime}=0$ and

$$
h_{j}\left(x_{1}, \ldots, x_{2 m}\right)=\sum_{i \leqslant m} \alpha_{i, j} x_{i}-\sum_{i \leqslant m} \alpha_{i, j} x_{2 m-i+1}
$$

for some $\alpha_{1, j}, \ldots, \alpha_{m, j} \in \mathbb{Q}$.
Proof. We know that each $h_{j}$ is a linear combination of $x_{1}, \ldots, x_{2 m}$. Essentially, Lemma 2.18 says that in $h_{j}\left(x_{1}, \ldots, x_{2 m}\right)$ the coefficients of $x_{i}$ and $x_{2 m-i+1}$ are the same except for their signs, which are different. Let $h_{j}\left(x_{1}, \ldots, x_{2 m}\right)=$ $\sum_{i \leqslant m} \alpha_{i} x_{i}-\sum_{i \leqslant m} \beta_{i} x_{2 m-i+1}$. Now let us fix some $i$. Consider the points $n \cdot\left(a_{1}, a_{2}, \ldots, a_{2 m}\right)$, where $a_{t}=1$ for $t=i, 2 m-i+1$ and $a_{t}=0$ otherwise. For $n \in \mathbb{Z}$ all
these points are elements of $V$ and $g_{j}$ holds for them. But $g_{j}$ takes for these points the form $n \cdot\left(\alpha_{i}-\beta_{i}\right)<c$ or $n \cdot\left(\alpha_{i}-\beta_{i}\right)=c^{\prime}$. Because it holds for every $n, \alpha_{i}$ and $\beta_{i}$ must be equal. In the last case we must have $c^{\prime}=0$.

It is a simple observation that if $g_{j}$ takes the form $h_{j}\left(x_{1}, \ldots, x_{2 m}\right)<c$, then $c>0$. Indeed, $g_{j}$ holds for $x_{1}=\cdots=x_{2 m}=0$ and $h_{j}(0, \ldots, 0)=0$.

Lemma 2.19. There are $m$ linearly independent polynomials among $h_{1}, h_{2}, \ldots, h_{\eta}$.
Proof. Consider the set $S=\left\{\boldsymbol{x} \in \mathbb{Q}^{2 m}: \forall i \leqslant \eta h_{i}(\boldsymbol{x})=0\right\}$. Clearly, $V \subseteq S$ and $S \subseteq V_{1}$. Assume that $\operatorname{dim} S>m$. Then there is a point $\boldsymbol{x} \in S \backslash V, \boldsymbol{x} \in X^{\prime}$. On the other hand, $V_{1} \cap X^{\prime} \subseteq V:$ a contradiction. So $\operatorname{dim} S=m$ and the number of linearly independent polynomials among $h_{1}, \ldots, h_{\eta}$ is equal to $2 m-\operatorname{dim} S=m$.

We may assume that $h_{1}, \ldots, h_{m}$ are linearly independent. Let condition $g_{i}$ be generated at the beginning of stage $s_{i}$. By the critical moment of $g_{i}$ we mean the beginning of stage $s_{i}$. We may assume that $s_{1}<s_{2}<\cdots<s_{m}$.

Definition 2.20. Suppose $H, H^{\prime}$ are heads of $M$. We say that a pair ( $H, H^{\prime}$ ) is dead at some moment of computation of $M$ if for some $l \leqslant 2 m$ head $H$ stands on the right side of block $B_{l}$ while head $H^{\prime}$ stands on the right side of block $B_{2 m-l+1}$.

Note that if a pair $\left(H, H^{\prime}\right)$ is dead at some moment, then it will remain dead for the rest of the computation. Intuitively, if a pair $\left(H, H^{\prime}\right)$ is dead then the heads $H$ and $H^{\prime}$ cannot be used to check that any two corresponding blocks are equal. However, it is not perfectly true. If variable $x_{j}$ occurs in $g_{j}$, then block $B_{j}$ must be read by some head before the critical moment of $g_{j}$. However, the length of $B_{j}$ may be recorded by some other head by its position in some other block. In turn, this information can be transmitted elsewhere and so on. Therefore, information about block $B_{j}$ can be used long after reading block $B_{j}$. This makes the analysis complex.

Now we shall show that $m \leqslant \frac{1}{2} k^{3}$. First consider $g_{1}$. Suppose that $x_{j}$ occurs in $g_{1}$ with a nonzero coefficient. Hence, there is a head $H$ which reads block $B_{j}$ before the critical moment of $g_{1}$. The coefficient of $x_{2 m-j+1}$ in $g_{1}$ is the same by Lemma 2.18, so also not equal to zero. Hence, before the critical moment of $g_{1}$ some head $H^{\prime}$ reads block $B_{2 m-j+1}$. We see that at the critical moment of $g_{1}$ the pair $\left(H, H^{\prime}\right)$ is dead. Take a look what happens next. We show that after at most $k$ next critical moments a new pair of heads becomes dead.

Assume that $H^{\prime \prime}$ is the head which reaches a new block at the beginning of stage $s_{1}$. Let $f_{1}, \ldots, f_{k-1}$ denote the distances of the other heads from the ends of the blocks they were in (some of $f_{i}$ 's might be equal to 0 ). If after this moment any reduced case condition is generated then it takes the form

$$
z+\sum \alpha_{i} f_{i}<c \quad \text { or } z+\sum \alpha_{i} f_{i}=c,
$$

where $\alpha_{i} \in \mathbb{Q}$ and $z$ is an expression depending only on the blocks read after the beginning of stage $s_{1}$. Consider some $g_{j}$ for $j>1$. Then $h_{j}=z+\sum \alpha_{i} f_{i}$, where $z$ is as above. Suppose $z \neq 0$. Then $z$ contains some variable $x_{t}$ standing for the length of block $B_{t}$ read by some head, say $H^{(3)}$, after reaching stage $s_{1}$. The corresponding block $B_{2 m-t+1}$ was read by some head $H^{(4)}$, not necessarily after reaching stage $s_{1}$, may be before. Note that if block $B_{2 m-t+1}$ was read by head $H^{(4)}$ after reaching stage $s_{1}$, then the pairs $\left(H^{(3)}, H^{(4)}\right)$ and $\left(H, H^{\prime}\right)$ must be different. Indeed, it is a consequence of the fact that the pair $\left(H, H^{\prime}\right)$ is dead after reaching stage $s_{1}$ and these heads cannot read any corresponding blocks.

Let $g_{t}$ be the first case condition such that at the critical moment of $g_{t}$ a new pair of heads becomes dead. Each $h_{j}$ for $1<j<t$ takes the form

$$
h_{j}=z_{j}+\sum_{i=1}^{k-1} \alpha_{i, j} f_{i},
$$

where $z_{j}$ contains only variables denoting the lengths of the blocks read after the beginning of stage $s_{1}$. As we have noticed, $z_{j}$ cannot contain simultaneously variables $x_{i}$ and $x_{2 m-i+1}(i \leqslant m)$ since otherwise a new pair of heads would be dead at the critical moment of $g_{j}$.

Lemma 2.21. Functions $u_{j}=\sum_{i=1}^{k-1} \alpha_{i, j} f_{i}$ for $1<j<t$ are linearly independent.
Proof. $\Lambda$ ssume that these functions are linearly dependent, i.e.

$$
\gamma_{2} u_{2}+\gamma_{3} u_{3}+\cdots+\gamma_{t-1} u_{t-1}=0
$$

for some $\gamma_{2}, \ldots, \gamma_{t-1} \in \mathbb{Q}$, not all equal 0 . We show that $\gamma_{2} h_{2}+\cdots+\gamma_{t-1} h_{t-1}=0$. Consider variables $x_{s}$ and $x_{2 m-s+1}(s \leqslant 2 m)$. Only one of them, say $x_{s}$, can be used in expressions $z_{j}(1<j<t)$. Variables $x_{s}$ and $x_{2 m-s+1}$ stand in $u_{j}$ with some coefficients $\lambda_{j}$ and $\Lambda_{j}$. Then

$$
\sum_{j=2}^{k-1} \gamma_{j} \lambda_{j}
$$

is the coefficient of $x_{s}$ in $\gamma_{2} u_{2}+\gamma_{3} u_{3}+\cdots+\gamma_{t-1} u_{t-1}$. So $\sum \gamma_{j} \lambda_{j}=0$. Similarly, $\sum \gamma_{j} \Lambda_{j}=0$. Note that now the coefficient of $x_{2 m-s+1}$ in $h_{j}$ is the same as in $u_{j}$, i.e. $\Lambda_{j}$, since $h_{j}=z_{j}+u_{j}$ and $x_{2 m-s+1}$ does not occur in $z_{j}$. The coefficient of $x_{s}$ in $h_{j}$ is $-\Lambda_{j}$ by Lemma 2.18. But $h_{j}=z_{j}+u_{j}$ and in $u_{j}$ variable $x_{s}$ has coefficient $\lambda_{j}$. So the coefficient of $x_{s}$ in $z_{j}$ is equal to $-\Lambda_{j}-\lambda_{j}$. On the other hand, we have

$$
\begin{aligned}
\gamma_{2} h_{2}+\cdots+\gamma_{t-1} h_{t-1} & =\gamma_{2}\left(z_{2}+u_{2}\right)+\cdots+\gamma_{t-1}\left(z_{t-1}+u_{t-1}\right) \\
& =\left(\gamma_{2} z_{2}+\cdots+\gamma_{t-1} z_{t-1}\right)+\left(\gamma_{2} u_{2}+\cdots+\gamma_{t-1} u_{t-1}\right) \\
& =\gamma_{2} z_{2}+\cdots+\gamma_{t-1} z_{t-1} .
\end{aligned}
$$

The last expression does not contain $x_{2 m-s+1}$ since no $z_{j}$ contains $x_{2 m-s+1}$. In turn the coefficient of $x_{s}$ is there equal to $\sum \gamma_{j}\left(-\Lambda_{j}-\lambda_{j}\right)=-\sum \gamma_{j} \Lambda_{j}-\sum \gamma_{j} \lambda_{j}=0$. Hence,
$\gamma_{2} h_{2}+\cdots+\gamma_{t-1} h_{t-1}$ does not contain $x_{s}$ and $x_{2 m-s+1}$. Number $s$ was arbitrary, so $\gamma_{2} h_{2}+\cdots+\gamma_{t-1} h_{t-1}=0$ and $h_{2}, \ldots, h_{t-1}$ are linearly dependent: a contradiction.

It follows from Lemma 2.21 that we have $t-2$ linearly independent functions $u_{j}$. Each $u_{j}$ is a linear combination of $f_{1}, \ldots, f_{k-1}$. Hence, $t-2 \leqslant k-1$, so $t \leqslant k+1$.
We have just proved that at the critical moment of $g_{k+1}$ a new pair of heads, different from ( $H, H^{\prime}$ ) must be dead. The above proof can be repeated virtually without change to show that for each $n$ at least one pair of heads becomes dead after the critical moment of $g_{n}$ and no later than the critical moment of $g_{n+k}$. There are only $k^{2} / 2$ pairs of heads (we allow the first and the second elements in a pair to be the same). Hence, $m$, the number of critical moments is not greater than $k \cdot k^{2} / 2=k^{3} / 2$. It completes the proof of Theorem 2.2.

## 3. Recognizing $P_{m}$

Recall that $P_{m, c}$ (for $c \in \mathbb{N}$ ) is a sublanguage of language $P_{m}$ defined as follows:

$$
P_{m, c}=\left\{1^{a_{1}} * \cdots * 1^{a_{2 m}} \in P_{m}: \forall n, l \leqslant 2 m a_{n}<c \cdot a_{l}\right\} .
$$

We show in this section that for each $c$ we can recognize language $P_{m, c}$ using an automaton with about $\sqrt[3]{m}$ heads. It shows that in some sense the bound given in Theorem 2.2 is stringent. However, $P_{m, c}$ is only a sublanguage of $P_{m}$ and we do not know such an algorithm for $P_{m}$.

Theorem 3.1. Let $k, c \in \mathbb{N}$. There is a $k$-head $1-M F A M$ recognizing a language $P^{\prime}$ such that $P_{m, c} \subseteq P^{\prime} \subseteq P_{m}$ and $m \geqslant k^{3} / 24$.

Proof. We consider only inputs of the form $1^{a_{1}} * \cdots * 1^{a_{2 m}}$, where $a_{l}<c \cdot a_{n}$ for each $n, l \leqslant 2 m$. By blocks of such an input word we mean the subwords $1^{a_{1}}, 1^{a_{2}}, \ldots, 1^{a_{2 m}}$. These blocks will be denoted by $B_{1}, B_{2}, \ldots, B_{2 m}$. Let $\left|B_{i}\right|$ stand for the length of block $B_{i}$. So $M$ has to check that for each $i \leqslant m,\left|B_{i}\right|=\left|B_{2 m-i+1}\right|$.
First we describe how $M$ can remember the length of some block $B_{l}$. For that purpose we need three heads, say $H_{1}, H_{2}, H_{3}$, with $H_{1}$ positioned at the beginning of $B_{l}$ and two other heads placed at the beginning of some block $B_{n}$. Information about the length of $B_{l}$ will be stored inside $B_{n}$. More precisely, we record only $\left[\left|B_{l}\right| / c\right]$ and $r\left(\left|B_{l}\right|, c\right)$ is to be remembered by the finite memory of $M$. Firstly, $H_{1}$ and $H_{2}$ move simultaneously until $H_{1}$ reaches the end of block $B_{l}$. For cach $c$ moves of $H_{1}$, head $H_{2}$ makes only one move to the right. So when $H_{1}$ reaches the end of $B_{1}$ then the distance between $H_{2}$ and $H_{3}$ is $\left[\left|B_{l}\right| / c\right]$. Head $H_{2}$ is still inside block $B_{n}$ since $\left|B_{l}\right|<c \cdot\left|B_{n}\right|$. Now we move simultaneously heads $H_{2}$ and $H_{3}$ with the same speed until $H_{2}$ reaches the end of $B_{n}$. The distance between $H_{2}$ and $H_{3}$ remains unchanged, so finally, $H_{3}$ is placed at the distance $\left[\left|B_{l}\right| / c\right]$ from the end of block $B_{n}$. If later we have to check that $\left|B_{l}\right|=\left|B_{2 m-l+1}\right|$, then we can use head $H_{3}$ and some other head $H$ placed at the
beginning of $B_{2 m-l+1}$. We start both of them, $H$ making $c$ moves for each single move of $H_{3}$. If $\left|B_{l}\right|=\left|B_{2 m-l+1}\right|$ then there are exactly $r\left(\left|B_{l}\right|, c\right)$ symbols left inside $B_{2 m-l+1}$ in front of head $H$ at the moment when $H_{3}$ reaches the end of $B_{i}$. This can be easily verified using the finite memory of $M$.

If $W$ is one of the blocks $B_{1}, B_{2}, \ldots, B_{2 m}$, say $B_{l}$, then let $W^{\prime}$ be the corresponding block $B_{2 m-l+1}$. The computation of $M$ consists of $k-1$ different stages. For each stage of execution there is a corresponding group of blocks. If $W$ belongs to such a group for stage $j$ then during stage $j$ it is checked whether $|W|=\left|W^{\prime}\right|$. Let the blocks corresponding to stage $j$ be denoted by $B_{j, 1}, B_{j, 2}, \ldots, B_{j, \sigma(j)}$. Each input word which we consider therefore takes the form

$$
\begin{aligned}
& \underbrace{B_{1,1} * B_{1,2} * \cdots * B_{1, \sigma(1)} *}_{\text {for stage } 1} * \underbrace{B_{2,1} * \cdots * B_{2, \sigma(2)} * \cdots * \underbrace{B_{k-1,1} * \cdots * B_{k-1, \sigma(k-1)}}_{\text {for stage } k-1} *}_{\text {for stage } 2} * \\
& \underbrace{B_{k-1, \sigma(k-1)}^{\prime} * \cdots * B_{k-1,1}^{\prime}}_{\text {for stage } k-1} * \underbrace{B_{k-2, \sigma(k-2)}^{\prime} * \cdots * B_{k-2,1}^{\prime}}_{\text {for stage } k-2} \cdots \cdots * \underbrace{B_{1, \sigma(1)}^{\prime} * \cdots * B_{1,1}^{\prime}}_{\text {for stage } 1}
\end{aligned}
$$

Values $\sigma(j)$ shall be determined later.
Now we describe one stage of computation, say stage $j$. It has the following important properties.
(i) Only heads $H_{1}, \ldots, H_{k-j+1}$ are in use, all of them initially placed at the beginning of block $B_{j, 1}$.
(ii) When stage $j$ ends, head $H_{k-j+1}$ is at the end of $B_{j, 1}^{\prime}$ and will not be used during the next stages, the remaining heads $H_{1}, \ldots, H_{k-j}$ are moved to the beginning of $B_{j+1,1}$.
For the sake of simplicity put $B_{j, i}=D_{i}$. Also let $t=k-j+1$. Stage $j$ is divided into $[(t-2) / 2]$ substages plus one additional "final" substage.


First substage: During this substage $M$ checks if the lengths of blocks $D_{\sigma(j)}, D_{\sigma(j)-1}$, $\ldots, D_{\sigma(j)-[t \mathrm{t}-2) / 2]+1}$ match with the lengths of the corresponding blocks. At the beginning, head $H_{t}$ moves to the end of $D_{\sigma(j)}$. In the meantime it reads blocks

$$
D_{\sigma(j)-[(t-2) / 2]+1}, \ldots, D_{\sigma(j)}
$$

and uses heads $H_{2}, H_{3}, \ldots, H_{t-1}$ to record their lengths inside $D_{1}$. It is possible since $2[(t-2) / 2]+1 \leqslant t-1$. Head $H_{1}$ is left unmoved at the beginning of $D_{1}$ and will stay there until the final substage. Then head $H_{t}$ reads blocks

$$
D_{\sigma(j)}^{\prime}, D_{\sigma(j)-1}^{\prime}, \ldots, D_{\sigma(i)-[(t-2) / 2]+1}^{\prime}
$$

and simultaneously, using information stored by the heads lying inside $D_{1}$ their lengths are checked.
Second substage: The second substage looks like the first one except for few details. At the beginning, head $H_{t-1}$ moves to the end of $D_{\sigma(j)}$. In the meantime it reads blocks

$$
D_{\sigma(j)-[(t-2) / 2]-[(t-4) / 2]+1}, \ldots, D_{\sigma(j)-\lfloor(t-2) / 2]}
$$

and uses heads $H_{3}, H_{4}, \ldots, H_{t-2}$ to record their length inside $D_{2}$. Again, it is possible since $2[(t-4) / 2]+2 \leqslant t-2$. Head $H_{2}$ is left unmoved at the beginning of $D_{2}$ and will stay there until the final substage. Then head $H_{t}$ reads blocks

$$
D_{\sigma(j)-[(t-2) / 2]}^{\prime}, \ldots, D_{\sigma(j)-[(t-2) / 2]-[(t-4) / 2]+1}^{\prime}
$$

and simultaneously, using information stored by the heads lying inside $D_{2}$, automaton $M$ checks the block lengths.

During subsequent substages $M$ works similarly. Each time one head is moved to the end of $D_{\sigma(j)}$ and one is left unmoved for the final substage. Also head $H_{t}$ reads some number of blocks. The remaining heads move one block forward because of the length checking. The number of such heads decreases by two each substage. Note that there must be at least four of them at the beginning of a substage. It follows that there are $[(t-2) / 2]$ of these substages.

Final substage: After the last nonfinal substage there are heads left at the beginnings of $D_{1}, D_{2} \ldots, D_{[(t-2) / 2]}$. Also in front of block $D_{[(t-2) / 2]+1}$ there are at least two heads. We move one of them to the beginning of $D_{[t / 2]+1}$. During this substage $M$ checks that the blocks $D_{1}, D_{2}, \ldots, D_{[t / 2]+1}$ have the same length as the corresponding blocks. It can be easily done since the number $\sigma(j)$ is chosen so that after the last nonfinal substage the head $H_{t}$ stands in front of $D_{[t / 2]+1}^{\prime}$. After the checking is done $M$ moves all heads (except $H_{t}$ and the other heads not already in use) to the end of $D_{\sigma(j)}$.

Now we count for how many blocks $D_{i}$ automaton $M$ checks that $\left|D_{i}\right|=\left|D_{i}^{\prime}\right|$ during stage $j$. During the final substage $M$ checks $[t / 2]+1$ blocks. During the first substage $M$ checks $[(t-2) / 2]$ blocks, during the second one only $[(t-4) / 2]$ of them, then $[(t-6) / 2],[(t-8) / 2], \ldots$ So during stage $j$ automaton $M$ checks together

$$
\sigma(j)=(1+[t / 2])+([(t-2) / 2]+[(t-4) / 2]+\cdots+1)
$$

blocks. Then

$$
\begin{aligned}
\sigma(j) & =([t / 2]+[(t-2) / 2]+[(t-4) / 2]+\cdots+1)+1 \\
& =\frac{1}{2} \cdot[t / 2] \cdot([t / 2]+1)+1 \\
& \geqslant \frac{1}{2} \cdot \frac{t-1}{2} \cdot\left(\frac{t-1}{2}+1\right)+1 \\
& =\frac{(t-1)(t+1)}{8}+1 \\
& =\frac{t^{2}-1}{8}+1 \geqslant \frac{t^{2}}{8}=\frac{(k-j+1)^{2}}{8} .
\end{aligned}
$$

The number of pairs of blocks checked during all stages can be estimated as follows:

$$
\begin{aligned}
\sum_{j=1}^{k-1} \sigma(j) & \geqslant \frac{k^{2}}{8}+\frac{(k-1)^{2}}{8}+\cdots+\frac{2^{2}}{8} \\
& =\frac{1}{8}\left(k^{2}+(k-1)^{2}+\cdots+1^{2}\right)-\frac{1}{8} \\
& =\frac{1}{8 \cdot 6} k \cdot(k+1) \cdot(2 k+1)-\frac{1}{8} \geqslant \frac{k^{3}}{24}
\end{aligned}
$$

The algorithm presented works effectively for large $k$. It can be improved slightly by combining it with a straightforward algorithm used for recognizing languages $L_{m}([16])$.

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[^0]:    * This research received support from the Institute of Informatics, University of Warsaw under program RP. I. 09 . The final version of the paper was prepared under sponsorship of the Alexander von HumboldtStiftung during author's stay at Technische Hochschule Darmstadt, Germany.

