
AUTOMATIC FUNCTIONS, LINEAR TIME AND LEARNING *

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ABSTRACT. The present work determines the exact nature of *linear time computable* notions which characterise automatic functions (those whose graphs are recognised by a finite automaton). The paper also determines which type of linear time notions permit full learnability for learning in the limit of automatic classes (families of languages which are uniformly recognised by a finite automaton). In particular it is shown that a function is automatic iff there is a one-tape Turing machine with a left end which computes the function in linear time where the input before the computation and the output after the computation both start at the left end. It is known that learners realised as automatic update functions are restrictive for learning. In the present work it is shown that one can overcome the problem by providing work tapes additional to a resource-bounded base tape while keeping the update-time to be linear in the length of the largest datum seen so far. In this model, one additional such work tape provides additional learning power over the automatic learner model and two additional work tapes give full learning power. Furthermore, one can also consider additional queues or additional stacks in place of additional work tapes and for these devices, one queue or two stacks are sufficient for full learning power while one stack is insufficient.

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1. INTRODUCTION

In inductive inference, automatic learners and linear time learners have played an important role, as both are considered as valid notions to model severely resource-bounded learners. On one hand, Pitt [28] observed that recursive learners can be made to be linear time learners by delaying; on the other hand, when learners are formalised by using automata updating a memory in each cycle with an automatic function, the corresponding learners are not as powerful as non-automatic learners [18] and cannot overcome their weakness by delaying. The relation between these two models is that automatic learners are indeed linear time learners [6] but not vice versa. This motivates to study the connection between linear time and automaticity on a deeper level.

It is well known that a finite automaton recognises a regular language in linear time. One can generalise the notion of automaticity from sets to relations and functions [3, 4, 15, 16, 21, 29] and say that a relation or a function is automatic iff an automaton recognises its graph, that is, if it reads all inputs and outputs at the same speed and accepts iff the inputs and outputs are related with each other, see Section 2 for a precise definition using the notion of convolution. For automatic functions it is not directly clear that they are in deterministic linear time, as recognising a graph and *computing the output of a string from the input* are two different tasks. Interestingly, in Section 2 below, it is shown that automatic functions coincide with those computed by linear time one-tape Turing machines which have the input and output both starting at the left end of the tape. In other words, a function is automatic iff it is linear-time computable with respect to the most restrictive variant of this notion; increasing the number of tapes or not restricting the position of the output on the tape results in a larger complexity class.

Section 3 is dedicated to the question on how powerful a linear time notion must be in order to capture full learning power in inductive inference. For the reader's convenience, here a short sketch of the underlying learning model is given: Suppose $L \subseteq \Sigma^*$ is a language. The learner gets as input a sequence x_0, x_1, \dots , of strings, where each string in L appears in the sequence and all the strings in the sequence are from L (such a sequence is called a text for L). As the learner is getting the input strings, it conjectures a sequence of grammars e_0, e_1, \dots as its hypotheses about what the input language is. These grammars correspond to some hypothesis space $\{H_e : e \in I\}$, where I is the set of possible indices and every possible learning task equals to some H_e . If this sequence of hypotheses converges to an index e for the language L (that is $H_e = L$), then one can say that the learner has learnt the input language from the given text. The learner learns a language L if it learns it from all texts for L . The learner learns a class \mathcal{L} of languages if it learns all languages from \mathcal{L} . The above is essentially the model of learning in the limit proposed by Gold [11]. Equivalently, one can consider the learner as operating in cycles, in n -th cycle it gets the datum x_n and conjectures the hypothesis e_n . In between the cycles, the learner may remember its previous inputs/work via some memory. The complexity of learners can be measured in terms of the complexity of mapping the old memory and input datum to the new memory and hypotheses. For automatic learners, one considers the above mapping to be given by an automatic function.

In respect to the automatic learners [6, 18, 19], it has been the practice to study the learnability of automatic classes (which are the set of all languages in some automatic family) and furthermore only to permit hypothesis spaces which are themselves automatic families containing the automatic class to be learnt. It turned out that certain automatic

families which are learnable by a recursive learner cannot be learnt by an automatic learner. The main weakness of an automatic learner is that it fails to memorise all past data. If one considers learning from fat text in which each datum occurs infinitely often, then automatic learners have the same learning power as recursive learners and their long-term memory can even be restricted to the size of the longest datum seen so far [18], the so called word size memory limitation.

Following the results of Section 2, one can simulate automatic learners by a learner using a one-tape Turing machine which updates the content of the tape in linear time in each round. In the present work this tape (called base tape) is restricted in length by the length of the longest datum seen so far — as the corresponding word size memory limitation of automatic learners studied in [18]. In each cycle, the learner reads one datum about the set to be learnt and revises its memory and conjecture. The question considered is how much extra power needs to be added to the learner for achieving full learnability; here the extra power is formalised by permitting additional work tapes which do not have length-restrictions; in each learning cycle the learner can, however, only work on these tapes in time linear in the length of the longest example seen so far. It can be shown using an archivation technique, that two additional work tapes can store all the data observed in a way that any learner can be simulated. When having only one additional work tape, the current results are partial: using a super-linear time-bound, one can simulate any learner for a class consisting entirely of infinite languages; furthermore, some classes not learnable by an automatic learner can be learnt using one work tape. When considering additional stacks in place of work tapes, two stacks are sufficient while one stack gives some extra learning power beyond that of an automatic learner but is insufficient to learn all in principle learnable classes.

2. AUTOMATIC FUNCTIONS AND LINEAR TIME

In the following, two concepts will be related to each other: automatic functions and functions computed by position-faithful one-tape Turing machines. In the following, a formal definition of these two concepts is given. Automatic functions and structures date back to the work of Hodgson [15, 16] and are based on the concept of convolution. A convolution permits to write pairs and tuples of strings by combining the symbols at the same position to new symbols.

Definition 2.1. Let Σ be a finite alphabet. Let \square be a special symbol not in Σ . The *convolution* $\text{conv}(x, y)$ of two strings $x = x_1x_2 \dots x_m$ and $y = y_1y_2 \dots y_n$, where $x_i, y_i \in \Sigma$, is defined as follows. Let $k = \max\{m, n\}$. For $i \in \{1, 2, \dots, k\}$, if $i \leq m$ then let $x'_i = x_i$ else let $x'_i = \square$; if $i \leq n$ then let $y'_i = y_i$ else let $y'_i = \square$. Now, the convolution is $\text{conv}(x, y) = (x'_1, y'_1)(x'_2, y'_2) \dots (x'_k, y'_k)$ where the symbols of this word are from the alphabet $(\Sigma \cup \{\square\}) \times (\Sigma \cup \{\square\})$.

Note that in the above definition, both x and y can be ε , the empty string. Similarly one can define the convolution of a fixed number of strings. Now the convolution permits to introduce automatic functions and relations.

Definition 2.2 (Hodgson [15, 16]). A function f , mapping strings to strings (possibly over a different alphabet), is said to be *automatic* iff the set $\{\text{conv}(x, f(x)) : x \in \text{dom}(f)\}$ is regular.

Similarly, an n -ary relation $R \subseteq \{(x_1, x_2, \dots, x_n) : x_1, x_2, \dots, x_n \in \Sigma^*\}$ is automatic iff $\{\text{conv}(x_1, x_2, \dots, x_n) : (x_1, x_2, \dots, x_n) \in R\}$ is regular. An n -ary function f is automatic iff $\{\text{conv}(x_1, x_2, \dots, x_n, f(x_1, x_2, \dots, x_n)) : (x_1, x_2, \dots, x_n) \in \text{dom}(f)\}$ is regular.

Here a regular set [17] is a set which is recognised by a deterministic finite automaton. This concept is equivalent to the one of sets recognised by non-deterministic finite automata. Furthermore, one can define regular sets inductively: Every finite set of strings is regular. The concatenation of two regular languages is regular, where $L \cdot H = \{xy : x \in L \wedge y \in H\}$; similarly, the union, intersection and set difference of two regular sets is regular. A further construct is the Kleene star, L^* , of a regular language L where $L^* = \{\varepsilon\} \cup L \cup L \cdot L \cup L \cdot L \cdot L \cup \dots = \{x_1 x_2 \dots x_n : x_1, x_2, \dots, x_n \in L\}$. Note that L^* always contains the empty string ε . The above mentioned operations are all that are needed, that is, a set is regular iff it can be constructed from finite sets by using the above mentioned operations in finitely many steps. The above inductive definition can be used to denote regular sets by regular expressions which are written representations of the above mentioned operations, for example, Σ^* represents the set of all strings over Σ and $\{00, 01, 10, 11\}^* \cap (\{0, 1\}^* \cdot \{0\} \cdot \{0, 1\}^*)$ represents the set of all binary strings of even length which contain at least one 0.

The importance of the concept of automatic functions and automatic relations is that every function or relation, which is first-order definable from a finite number of automatic functions and relations, is automatic again and the corresponding automaton can be computed effectively from the other automata. This gives the second nice fact that structures consisting of automatic functions and relations have a decidable first-order theory [16, 21].

A *position-faithful one-tape Turing machine* is a Turing machine which uses a one-side infinite tape, with the left-end having a special symbol \boxplus which only occurs at this position and cannot be modified. The input starts from the cell at the right of \boxplus and is during the computation replaced by the output which starts from the same cell. The end of input and output is the first appearance of the symbol \boxminus which is the default value of an empty cell before it is touched by the head of the Turing machine or filled with the input.

It is assumed that the Turing machine halts when it enters an accepting/final state (if ever). A position-faithful one-tape Turing machine computes a function f , if when started with tape content being $\boxplus x \boxminus^\infty$, the head initially being at \boxplus , the Turing machine eventually reaches an accepting state (and halts), with the tape content starting with $\boxplus f(x) \boxminus$. Note that there is no restriction on the output beyond the first appearance of \boxminus . Furthermore, a Turing machine can halt without reaching an accepting state, in which case the computation is not valid; this possibility is needed when a non-deterministic Turing machine has to satisfy a time bound on the duration of the computation.

Though the requirement of “position-faithfulness” seems to be a bit artificial, it turns out that it is a necessary requirement. This is not surprising, as moving i bits by j cells requires, in the worst case, proportional to $i \cdot j$ steps. So sacrificing the requirement of position-faithfulness clearly increases the class. For example, the function which outputs the binary symbols between the first and second occurrence of a digit other than 0 and 1 of an input would become linear time computable by an ordinary one-tape Turing machine although this function is not linear time computable by a position-faithful one-tape Turing machine. Such an additional side-way to move information (which cannot be done in linear time on one-tape machines) has therefore been excluded from the model. The functions computed by position-faithful one-tape Turing machines are in a certain sense a small natural class of linear time computable functions.

Some examples of automatic functions are those which append to or delete in a string

some characters as long as the number of these characters is bounded by a constant. For example a function deleting the first occurrence (if any) of 0 in a string would be automatic; however, a function deleting all occurrences of 0 is not automatic. Below is a more comprehensive example of an automatic function.

Example 2.3. Suppose $\Sigma = \{0, 1, 2\}$. Suppose f is a mapping from Σ^* to Σ^* such that $f(x)$ interchanges the first and last symbol in x ; $f(\varepsilon) = \varepsilon$. Then f is automatic and furthermore f can also be computed by a position-faithful one-tape Turing machine.

To see the first, note that the union of the set $\{\varepsilon\} \cup \{(a, a) : a \in \Sigma\}$ and all sets of the form $\{(a, b)\} \cdot \{(0, 0), (1, 1), (2, 2)\}^* \cdot \{(b, a)\}$ with $a, b \in \Sigma$ is a regular set. Thus $\{\text{conv}(x, y) : x \in \Sigma^* \wedge y = f(x)\}$ is a regular set and f is automatic.

A position-faithful one-tape Turing machine would start on the starting symbol \boxplus and go one step right. In the case that there is a \boxminus in that cell, the machine halts. Otherwise it memorises in its state the symbol a there. Then it goes right until it finds \boxminus ; it then goes one step left. The Turing machine then memorises the symbol b at this position and replaces it by a . It then goes left until it finds \boxplus , goes one step right and writes b .

That f in the preceding example can be computed in both ways is not surprising, but indeed a consequence of the main result of this section which states that the following three models are equivalent:

- automatic functions;
- functions computed in deterministic linear time by a position-faithful one-tape Turing machine;
- functions computed in non-deterministic linear time by a position-faithful one-tape Turing machine.

This equivalence is shown in the following two results, where the first one generalises prior work [6, Remark 2].

Theorem 2.4. *Let f be an automatic function. Then there is a deterministic linear time one-tape position-faithful Turing machine which computes f .*

Proof. The idea of the proof is to simulate the behaviour of a deterministic finite automaton recognising the graph of f . The Turing Machine goes from the left to the right over the input word and takes note of which states of the automaton can be reached from the input with only one unique possible output. Once the automaton reaches an accepting state in this simulation (for input/output pairs), the simulating Turing machine turns back (that is, it goes from right to left over the tape) converting the sequence of inputs and the stored information about states as above into that output which produces the unique accepting run on the input. Now the formal proof is given.

Suppose that a deterministic automaton with c states (numbered 1 to c , where 1 is the starting state) accepts a word of the form $\text{conv}(x, y) \cdot (\boxminus, \boxminus)$ iff x is in the domain of f and $y = f(x)$; the automaton rejects any other sequence. Note that this small modification of the way the convolution is represented simplifies the proof. As $f(x)$ depends uniquely on x , any string of the form $\text{conv}(x, y) \cdot (\boxminus, \boxminus)$ accepted by the automaton satisfies $|y| \leq |x| + c$. Let δ be the transition function for the automaton above and $\hat{\delta}$ be the corresponding extended transition function [17].

Suppose that the input is $x = x_1x_2 \dots x_r$. Let the cell number k be that cell which carries the input x_k (with cell 0 carrying \boxplus), that is the k -th cell to the right of \boxplus ; \boxplus is in cell number 0. Note that the Turing Machine described below does not use the cell number

in its computation; the numbering is used just for ease of notation. The simulating Turing machine uses a larger tape alphabet containing extra symbols from $(\Sigma \cup \square) \times \{+, -, *\}^c$, that is, one considers the additional symbols consisting of tuples of the form $(a, s_1, s_2, \dots, s_c)$, where $a \in \Sigma \cup \{\square\}$ and $s_i \in \{-, +, *\}$. These symbols are written temporarily onto the tape while processing the word from the left to the right and later replaced when coming back from the right to the left.

Intuitively, during the computation while going from left to right, for cell number k , one wishes to replace x_k by the tuple $(x_k, s_1^k, s_2^k, \dots, s_c^k)$ where, for $d \in \{1, 2, \dots, c\}$: $s_d^k = -$ iff there is no word of the form $y_1 y_2 \dots y_{k-1}$ such that the automaton on input $(x_1, y_1)(x_2, y_2) \dots (x_{k-1}, y_{k-1})$ reaches the state d (that is, for no $y_1 y_2 \dots y_{k-1}$, $\hat{\delta}(1, (x_1, y_1)(x_2, y_2) \dots (x_{k-1}, y_{k-1})) = d$); $s_d^k = +$ iff there is exactly one such word; $s_d^k = *$ iff there are at least two such words. Here the x_i and y_i can also be \square (when i is larger than the length of the relevant string, for example $x_i = \square$ for $i > r$).

For doing the above, the Turing machine simulating the automaton replaces the cell to the right of \boxplus , that is, the cell containing x_1 , by $(x_1, +, -, \dots, -)$. Then, for the k -th cell, $k > 1$, to the right of \boxplus , with entry x_k (from the input or \square if $k > r$) the Turing machine replaces $x s_k$ by $(x_k, s_1^k, s_2^k, \dots, s_c^k)$ under the following conditions, (where the entry in the cell to the left was $(x_{k-1}, s_1^{k-1}, s_2^{k-1}, \dots, s_c^{k-1})$ and where d ranges over $\{1, 2, \dots, c\}$):

- s_d^k is $+$ iff there is exactly one $(y_{k-1}, d') \in (\Sigma \cup \{\square\}) \times \{1, 2, \dots, c\}$ such that $s_{d'}^{k-1}$ is $+$ and $\delta(d', (x_{k-1}, y_{k-1})) = d$ and there is no pair $(y_{k-1}, d') \in (\Sigma \cup \{\square\}) \times \{1, 2, \dots, c\}$ such that $s_{d'}^{k-1}$ is $*$ and $\delta(d', (x_{k-1}, y_{k-1})) = d$;
- s_d^k is $*$ iff there are at least two pairs $(y_{k-1}, d') \in (\Sigma \cup \{\square\}) \times \{1, 2, \dots, c\}$ such that $s_{d'}^{k-1}$ is $+$ and $\delta(d', (x_{k-1}, y_{k-1})) = d$ or there is at least one pair $(y_{k-1}, d') \in (\Sigma \cup \{\square\}) \times \{1, 2, \dots, c\}$ such that $s_{d'}^{k-1}$ is $*$ and $\delta(d', (x_{k-1}, y_{k-1})) = d$;
- s_d^k is $-$ iff for all pairs $(y_{k-1}, d') \in (\Sigma \cup \{\square\}) \times \{1, 2, \dots, c\}$ such that $\delta(d', (x_{k-1}, y_{k-1})) = d$, it holds that $s_{d'}^{k-1}$ is $-$.

Note that the third case applies iff the first two do not apply. The automaton replaces each symbol in the input as above until it reaches the cell where the intended symbol $(x_k, s_1^k, s_2^k, \dots, s_c^k)$ has $s_d^k = +$ for some accepting state d . (Note that the accepting states occur in the automaton only if both the input and output are exhausted by the convention made above.) If this happens, the Turing machine turns around, memorises the state d , erases this cell (that is, writes \square) and goes left.

When the Turing machine moves left from the cell number $k + 1$ to the cell number k (which contains the entry $(x_k, s_1^k, s_2^k, \dots, s_c^k)$), where the state memorised for the cell number $k + 1$ is d' , then it determines the unique $(d, y_k) \in \{1, 2, \dots, c\} \times (\Sigma \cup \{\square\})$ such that $s_d^k = +$ and $\delta(d, (x_k, y_k)) = d'$; then the Turing machine replaces the symbol on cell k by y_k . Then the automaton keeps the state d in the memory and goes to the left and repeats this process until it reaches the cell which has the symbol \boxplus on it. Once the Turing machine reaches there, it terminates.

For the verification, note that the output $y = y_1 y_2 \dots$ (with \square appended) satisfies that the automaton, after reading $(x_1, y_1)(x_2, y_2) \dots (x_k, y_k)$, is always in a state d with $s_d^{k+1} = +$, as the function value y is unique in x ; thus, whenever the automaton ends up in an accepting state d with $s_d^{k+1} = +$ then the input-output-pair $\text{conv}(x, y) \cdot (\square, \square)$ has been completely processed and $x \in \text{dom}(f) \wedge f(x) = y$ has been verified. Therefore, the Turing machine can turn and follow the unique path, marked by $+$ symbols, backwards in order to reconstruct the output from the input and the markings. All superfluous symbols and markings are

removed from the tape in this process. As the automaton accepts $\text{conv}(x, y) \cdot (\square, \square)$, and y depends uniquely on x , $|y| \leq |x| + c$. Hence the runtime of the Turing machine is bounded by $2 \cdot (|x| + c + 2)$, that is, the runtime is linear. \square

Remark 2.5. Note that the Turing machine in the above theorem makes two passes, one from the origin to the end of the word (plus maybe constantly many more cells) and one back. These two passes are needed for a deterministic Turing machine: Recall the function f from Example 2.3 with $f(x_1x_2 \dots x_{k-1}x_k) = x_kx_2 \dots x_{k-1}x_1$ for all non-empty words $x_1x_2 \dots x_{k-1}x_k$. When starting at the left end, the machine has first to proceed to the right end to read the last symbol before it can come back to the left end in order to write that symbol into the new first position. Hence the runtime of the one-tape deterministic Turing machine (for the simulation as in Theorem 2.4) cannot be below $2 \cdot |x|$ for an input x . Non-deterministic Turing machines can, however, perform this task with one pass.

For the converse direction of the equivalence of the two models of computation, that is, of automatic functions and position-faithful linear time Turing machines, assume that a function is computed by a position-faithful non-deterministic one-tape Turing machine in linear time. For an input x , any two non-deterministic accepting computations have to produce the same output $f(x)$. Furthermore, the runtime of each computation has to follow the same linear bound $c \cdot (|x| + 1)$, independent of whether the computation ends up in an accepting state or a rejecting state.

Theorem 2.6. *Let f be a function computed by a non-deterministic one-tape position-faithful Turing machine in linear time. Then f is automatic.*

Proof. The proof is based on crossing-sequence methods, see [13, 14] and [26, Section VIII.1]. The idea is to show that f is automatic by providing a non-deterministic automaton which recognises the graph f of the function by going symbol by symbol over the convolution of input and output and for each symbol, the automaton guesses, for the Turing Machine on the corresponding input, the crossing sequence on the right side and verifies that this crossing-sequence is compatible with the previously guessed crossing-sequence on the left side of the symbol plus the local transformation of the respective input symbol to the output symbol. This is now explained in more detail.

Without loss of generality one can assume that the position-faithful Turing machine M computing f starts at \boxplus and returns to that position at the end; a computation accepts only when the full computation has been accomplished and the automaton has returned to \boxplus . By a result of Hartmanis [12] and Trakhtenbrot [30], there is a constant c' such that an accepting computation visits each cell of the tape at most c' times; otherwise the function f would not be linear time computable. This permits to represent the computation locally by considering for each visit to a cell — the direction from which the Turing machine M entered the cell, in which state it was, what activity it did and in which direction it left the cell. Below, the k -th cell to the right of \boxplus is referred to as cell number k . The local computation at the cell number k can be considered as a tuple $(x_k, is_k^1, os_k^1, d_k^1, z_k^1, is_k^2, os_k^2, d_k^2, z_k^2, \dots, is_k^{r_k}, os_k^{r_k}, d_k^{r_k}, z_k^{r_k})$, for some $r_k \leq c'$, where x_k is the initial symbol at the cell number k , and for each j , is_k^j denotes the state the Turing machine M was in when it visited the cell number k for the j -th time, os_k^j is the state the Turing machine M was in when it left the cell number k after the j -th visit, d_k^j is the direction in which the Turing machine M left after the j -th visit, and z_k^j is the symbol written in the cell number k by the Turing machine M during the j -th visit; r_k here denotes the total number of visits of the Turing machine M to the k -th cell. Note that

the number of possibilities for the local computation as above is bounded by a constant.

As an intermediate step one shows that a non-deterministic finite state automaton can recognise the set

$$A = \{\text{conv}(x \cdot \square^s, y \cdot \square^{s+|x|-|y|}) : x \in \text{dom}(f) \wedge y = f(x) \wedge s > 0 \wedge s + |x| > |y| \wedge \text{the Turing machine } M \text{ on input } x \text{ does not move beyond cell number } s - 1 + |x|\}.$$

This is done by initially guessing the local computation at \boxplus (the 0-th cell). Then on each subsequent input (x_k, y_k) (where x_k or y_k might be \square), starting with $k = 1$, the automaton (i) guesses the local computation at the k -th cell, (ii) checks that this guess in (i) is consistent with the local computation guessed at cell $k - 1$ (that is, each time the Turing machine M moved from cell $k - 1$ to cell k or cell k to $k - 1$, the corresponding guessed leaving/entering states match), (iii) the computation within the cell is consistent with the Turing machine M 's state table (that is, either each of the entries $is_k^j, os_k^j, d_k^j, z_k^j$ satisfies that Turing machine has transition from state is_k^j on reading input z_k^{j-1} to state os_k^j writing z_k^j in the cell and moving in direction d_k^j , where $z_k^0 = x_k$ and $z_k^k = y_k$ or the Turing machine does not reach this cell and $y_k = x_k$), (iv) for the last input the automaton also checks that it is of the form (\square, \square) and that the Turing machine M does not reach this cell.

If at the end, all the computation and guesses are consistent then the automaton accepts. The automaton thus passes over the full word and accepts $\text{conv}(x \cdot \square^s, y \cdot \square^{s+|x|-|y|})$ iff the non-deterministic computation transforms $\boxplus x \square^s$ into $\boxplus y \square^{s+|x|-|y|}$.

It follows that the set $B = \{\text{conv}(x, y) : x \in \text{dom}(f) \wedge y = f(x)\}$ is regular as well, as it is first-order definable from A and the prefix relation: $z \in B \Leftrightarrow z$ does not end with (\square, \square) and $z \cdot (\square, \square)$ is a prefix of an element in A . Thus f is automatic. \square

Remark 2.7. One might ask whether the condition on the input and output starting at the same position is really needed. The answer is “yes”. Assume by way of contradiction that it would not be needed and that all functions linear time computable by a one-tape Turing machine without any restrictions on output positions are automatic. Then one could consider the free monoid over $\{0, 1\}$. For this monoid, the following function could be computed from $\text{conv}(x, y)$: The output is $z = f(x, y)$ if $y = xz$; the output is $\#$ if such a z does not exist. For this, the machine just compares x_1 with y_1 and erases (x_1, y_1) , x_2 with y_2 and erases (x_2, y_2) and so on, until it reaches (a) a pair of the form (x_m, y_m) with $x_m \neq y_m$ or (b) a pair of the form (x_m, \square) or (c) a pair of the form (\square, y_m) or (d) the end of the input. In cases (a) and (b) the output has to be $\#$ and the machine just erases all remaining input symbols and puts the special symbol $\#$ to denote the special case; in case (c) the value z is just obtained by changing all remaining input symbols $(\#, y_k)$ to y_k and the Turing machine terminates. In case (d) the valid output is the empty string and the Turing machine codes it adequately on the tape. Hence f would be automatic. But now one could first-order define concatenation g by letting $g(x, z)$ be the y for which $f(x, y) = z$; this would give that the concatenation is automatic, which is known to be false. The non-automaticity of the concatenation can be seen as follows: For each automatic function there is, by the pumping lemma [17], a constant c such that each value is at most c symbols longer than the corresponding input; now the mapping $\text{conv}(x, y) \mapsto xy$ fails to satisfy this for any given constant c , for example, $x = 0^{c+1}$ and $y = 1^{c+1}$ are mapped to $xy = 0^{c+1}1^{c+1}$. Hence the condition on the starting-positions cannot be dropped.

One can generalise non-deterministic computation to computation by alternating Turing machines [8]. Well known results in this field [8] are that sets decidable in alternating logarithmic space are equal to sets decidable in polynomial time and that alternating polynomial time computations define the class PSPACE for sets. Therefore it might be interesting to ask what is the equivalent notion for alternating linear time computation. The following definition deals with the alternating computation counterpart of position-faithful linear time computations.

Definition 2.8. An alternating position-faithful one-tape Turing machine M has \exists -states and \forall -states among the Turing machine states which permit the machine to guess one bit (which can then be taken into account in future computation). It uses, as the name says, exactly one tape which initially contains $\boxplus x \boxminus^\infty$, where x is the input string. At the end of the computation, the output is the string between the \boxplus and the first \boxminus . M is linear time bounded iff there is a constant c such that, for each input x of length n and each run of M , the duration of the run until M halts is at most $c \cdot (n + 1)$ time steps. Furthermore, M alternately computes a function f iff for each string x on the input there is a unique string y (which must be equal to $f(x)$) such that, for a computation tree T formed by choosing at each \exists -state the guessed bit appropriately (the \forall -states are still true branching nodes in this tree T), one has that each computation path on T ends up in an accepting state and each computation produces the same output y .

It is easy to see that every function f computed non-deterministically by a position-faithful one-tape Turing machine in linear time is also computed by an alternating position-faithful one-tape Turing machine in linear time. However, the converse direction is open; if the answer would be negative, one could use it as the basic definition of a concept similar to automatic structures which is slightly more general.

Open Problem 2.9. Is every function f computable in alternating linear time by a position-faithful one-tape Turing machine automatic?

3. LINEAR TIME LEARNERS

The following definition of learning is based on the Gold's [11] notion of learning in the limit. The presentation differs slightly in order to incorporate memory restrictions and automaticity as considered in this paper; note that learners without any restrictions on the way the long term memory is organised can store all past data and are therefore as powerful as those considered by Gold [11].

Informally, a learning scenario can be described as follows. Suppose a family $\{L_e : e \in I\}$ of languages is given (in some effective form), where I is an index set. The learner, as input, gets a listing of the elements of some set L_e . The learner is supposed to figure out, in the limit from the listing as above, an index d such that $L_d = L_e$. For ease of presentation it is assumed that all languages L_e are not empty.

The listing of elements is formalised as a text. A *text* T for a language L is an infinite sequence, w_0, w_1, w_2, \dots , containing all elements of L but no non-element of L , in any order with repetitions allowed. Let $T[n]$ denote the sequence of first n elements of the text: w_0, w_1, \dots, w_{n-1} . The basic model of inductive inference [1, 2, 11, 20, 27] is that the learner M is given a text w_0, w_1, \dots of all the words in a language L , one word per cycle. At the same time M outputs a sequence e_0, e_1, \dots of indices, one index in each cycle. Intuitively,

each e_i can be considered as a conjecture of the learner regarding what the language L is, based on the data w_0, w_1, \dots, w_{i-1} . In general, the indices conjectured are from some index set J and interpreted in a hypothesis space $\{H_e : e \in J\}$, where $\{L_e : e \in I\} \subseteq \{H_e : e \in J\}$.

The learner maintains information about past data in form of some memory, which may change between cycles. Thus, the learner can be considered as an

algorithmic mapping from (old memory, new datum) to (new memory, new conjecture)

where the learner has some fixed initial memory. The learner *learns* or *identifies* the language L , if, for all possible texts for L , there is some k such that e_k is an index for L and $e_{k'} = e_k$ for $k' \geq k$. The learner learns a class \mathcal{L} of languages if it learns each language in \mathcal{L} .

The most basic set of hypothesis spaces are automatic families of languages. Here, a family of languages $\{L_e : e \in I\}$, is *automatic* if the index set I and the set $\{\text{conv}(e, x) : e \in I, x \in L_e\}$ are both regular. Automatic families [18, 19] are the automata-theoretic counterpart of indexed families [1, 23] which were widely used in inductive inference to represent the class to be learnt. Note that when $\{H_d : d \in J\}$ is a hypothesis space for $\{L_e : e \in I\}$, which is an automatic family as well, then there is an automatic function f mapping the indices from J back to those in I , that is, $L_{f(d)} = H_d$ for all those $d \in J$ where H_d equals some L_e . Hence one can without loss of generality (for learning criteria considered in this paper) directly use the hypothesis space $\{L_e : e \in I\}$ for the class $\{L_e : e \in I\}$ to be learnt.

A learner M is called *automatic* if the mapping (old memory, new input word) to (new memory, new conjecture) for the learner is automatic, that is, the set

$$\{\text{conv}(om, dat, nm, nc) : M(om, dat) = (nm, nc)\}$$

is regular. In general, om and nm are the old and new versions of the long term memory of the learner. Automatic learners are, roughly speaking, the most restrictive form of learners which update a long term memory in each cycle where they process one new datum given to the learner.

The next definition generalises the notion of automatic learning to a learner which has a linear or nearly linear time bound for each of its cycle. This generalisation is natural, due to the correspondence between automatic function and linear time computable functions given in the previous section of this paper.

Definition 3.1. A learner M is a Turing machine which maintains some memory and in each cycle receives as input one word to be learnt, updates its memory and then outputs an hypothesis. The tapes of the Turing machine are all one-sided infinite and contain \boxplus at the left end. The machine operates in cycles, where in each cycle it reads one current datum (from a text of the language to be learnt) and formulates one hypothesis. Furthermore, it has some long term memory in its tape where the memory in Tape 0 is always there while the memories in the additional data structures (Tapes 1, 2, \dots , k) is only there when these additional data structures are explicitly permitted.

- At the beginning of each cycle, Tape 0 (base tape) contains convolution of the input and some information (previous long term memory) which is not longer in length (up to an additive constant) than the length of the longest word seen so far. The head on Tape 0 of the Turing machine starts at \boxplus at the beginning of each cycle.
- At the end of each cycle, Tape 0 (base tape) has to contain the convolution of the new long term memory and the hypothesis which the learner is conjecturing.

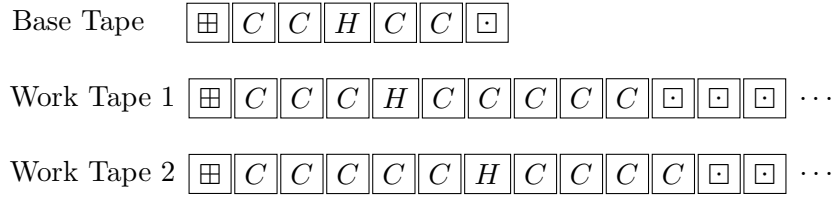


Figure 1: Learner with two working tapes; the head positions are H and other data positions are C ; note that data characters can be convoluted characters from finitely many alphabets in order to store the convolution of several tracks on a tape.

- During the execution of a cycle, the learner can run in time linear in the current length of Tape 0 and, like a position-faithful one-tape Turing machine, replace the convolution of the current datum and old long term memory by the convolution of the hypothesis and the new long term memory. Furthermore, the memory in Tape 0 has to meet the constraint that it is at most (up to an additive constant) the length of the longest datum seen so far (including the current datum), hence there is an explicit bound on the length of Tape 0 in each cycle.
- Tapes 1, 2, \dots , k are normal tapes, whose contents and head positions are not modified during change of cycles. M can use these tapes for archiving information and doing calculations. There is no extra time allowance for the machine to use these tapes, hence the machine can only access a small amount (linear in the length of Tape 0) in each cycle of these tapes.
- Without loss of generality, one can assume that the length of the longest datum seen so far is stored in the memory in Tape 0.

The learner is said to have k additional work tapes iff it has in addition to Tape 0 also the Tapes 1, 2, \dots , k . Figure 1 illustrates a learner with two additional tapes.

Note that in the definition of Tape 0, it is explicit that the length of the hypothesis produced is bounded by the length of the largest example seen so far plus a constant. This is compatible with learning, as for all automatic families, (i) for any finite set L in the family, the length of the smallest index e for L overshoots the length of the longest element of L by at most a constant and (ii) for any infinite set L in the family there are words in L which are longer than some index for L ; thus a learner cannot fail just because the indices of a language L_e are extremely long compared to the size of the members of L_e – though, of course, there may be other reasons for a learner not to be successful.

Note that if only Tape 0 is present, the model is equivalent to an automatic learner with the memory bounded by the size of the longest datum seen so far (plus a constant) [6, 18]. The next examples illustrate what type of learnable automatic classes exist.

Example 3.2. The following automatic classes are learnable by an automatic learner with its memory bounded by the length of the longest example seen so far (plus a constant):

First, the class of all extensions of an index, that is, $I = \Sigma^*$ and $L_e = e \cdot \Sigma^*$ for all $e \in I$. Here the learner maintains as a memory the longest common prefix e of all data seen so far and whenever the memory is e and a new datum x is processed, the learner updates e to the longest common prefix of both, e and x , which is also the next hypothesis.

Second, the class of all closed intervals in the lexicographic ordering, that is, $I = \{conv(d, e) : d, e \in \Sigma^* \wedge d \leq_{lex} e\}$ and $L_{conv(d, e)} = \{x \in \Sigma^* : d \leq_{lex} x \leq_{lex} e\}$; here

$x \leq_{lex} y$ denotes that x is lexicographically before y . The learner maintains as memory the lexicographically least and greatest elements seen so far, the convolution of these elements also serves as hypothesis.

Third, the class of all strings of length different from the index, that is, $I = \{0\}^*$ and $L_e = \{x \in \Sigma^* : |x| \neq |e|\}$. Here the learner archives in Tape 0 binary string which is of length 1 plus the length of the longest example seen so far; the k -th bit of this string (starting with $k = 0$) is 1 iff an example of length k has been seen so far, and 0 iff no example of length k has been seen so far. The conjecture is 0^h for the least h such that either the h -th bit of the memory is 0 or h is 1 plus the length of the memory string.

For any automatic family, $\{H_e : e \in J\}$, the equivalence question for indices is automatic, that is, the set $\{\text{conv}(e, e') : H_e = H_{e'}\}$ is regular. Thus for the purposes of this paper, one can take the hypothesis space to be one-one, that is, different indices represent different languages. In a one-one hypothesis space, the index e of a finite language L_e has, up to an additive constant, the same length as the longest word in L_e ; this follows easily from [19, Theorem 3.5]. This observation is crucial as otherwise the time-constraint on the learner would prevent the learner from eventually outputting the correct index; for infinite languages this is not a problem as the language must contain arbitrarily long words.

Angluin [1] gave a characterisation when a class is learnable in general. This characterisation, adjusted to automatic families, says that a class is learnable iff, for every $e \in I$, there exists a finite set $D \subseteq L_e$ such that there is no $d \in I$ with $D \subseteq L_d \subset L_e$. All the automatic families from Example 3.2 satisfy this criterion; however, Gold [11] provided a simple example of a non-learnable class which of course then also fails at Angluin's criterion: One infinite set plus all of its finite subsets.

The main question of this section is which learnable classes can also be learnt by a linear-time learner with k additional work tapes. For $k = 0$, this is in general not possible, as automatic learners fail to learn various learnable classes [18], for example the class of all sets $\{0, 1\}^* - \{x\}$, with the index x being from $\{0, 1\}^*$, and the class of all sets $L_e = \{x \in \{0, 1\}^{|e|} : x \neq e\}$.

Freivalds, Kinber and Smith [10] introduced limitations on the long term memory into inductive inference; Kinber and Stephan [22] transferred it to the field of language learning. Automatic learners have similar limitations and are therefore not able to learn all learnable automatic classes [6, 18]. The usage of additional work tapes for linear time learners permits to overcome these limitations, the next results specify how many additional work tapes are needed. Recall from above that work tapes are said to be *additional* iff they are in addition to the base tape.

Theorem 3.3. *Suppose $\Sigma = \{0, 1, 2\}$ and consider the automatic family \mathcal{L} over the alphabet Σ which is defined as follows: \mathcal{L} consists of (i) $L_\varepsilon = \{0, 1\}^*$ and (ii) $L_{x0} = \{0, 1\}^* \cup \{x2\} - \{x\}$ and (iii) $L_{x1} = \{0, 1\}^* \cup \{x2\}$, for each $x \in \{0, 1\}^*$. Then, \mathcal{L} does not have an automatic learner but has a linear-time learner using one additional work tape.*

Proof. An automatic learner cannot memorise all the data from $\{0, 1\}^*$ it sees. For any automatic learner, one can show, see [18], that there are two finite sequences of words from L_ε , one containing x and one not containing x , such that the automatic learner has the same long term memory after having seen both sequences. If one presents to the automatic learner, after these sequences, all the elements of L_{x0} , then the automatic learner's limiting behaviour on the two texts so formed is the same, even though they are texts for two

different languages, L_{x1} or L_{x0} , in \mathcal{L} . Therefore the learner cannot learn the class \mathcal{L} .

A linear time learner with one additional work tape (called Tape 1) initially conjectures L_ε and uses Tape 1 to archive all the examples seen at the current end of the written part of the tape. When the learner sees a word of the form $x2$, it maintains a copy of it in the memory part of Tape 0 and conjectures $x0$ as its hypothesis. In each subsequent cycle, the learner scrolls back Tape 1 by one word and compares the word there as well as the current input with $x2$; if one of these two is x then the learner changes its conjecture to L_{x1} , else it keeps its conjecture as L_{x0} . In the case that the origin of Tape 1 (\boxplus) is reached, the learner from then onwards ignores Tape 1 and only compares the incoming input with $x2$. It is easy to verify that the learner as described above learns \mathcal{L} . \square

Theorem 3.4. *Every learnable automatic family \mathcal{L} has a linear-time learner using two additional work tapes.*

Proof. Jain, Luo and Stephan [18] showed that for every learnable automatic family $\mathcal{L} = \{L_e : e \in I\}$ there is an automatic learner M using memory bounded in length by the length of the longest example seen so far (plus a constant) which learns the class from every fat text (a text in which every element of the language appears infinitely often). So the main idea is to use the two additional tapes in order to simulate and feed the learner M with a fat text. The two additional tapes are used to store all the incoming data and then to feed the learner M with each data item infinitely often. The words in the tapes are stored using some separator $\#$ to separate the words. Thus, $00\#1\#\#11\boxplus$ indicates that the tape contains the words 00 , 1 , ε and 11 .

The learner N for \mathcal{L} using two additional tapes works as follows. Suppose the previous memory stored in Tape 0 is mem_k (initially the memory stored on Tape 0 is the initial memory of M) and the current datum is w_k . Then, N does the following:

- Compute $M(mem_k, w_k) = (mem', e')$.
- Find the last word in Tape 1, say t . Erase this word from Tape 1. In the case that Tape 1 was already empty, let $t = w_k$.
- Compute $M(mem', t) = (mem_{k+1}, e_k)$.
- Write w_k and t at the end of Tape 2 (using the separator $\#$ to separate the words).
- When the beginning of Tape 1 is reached (\boxplus), interchange the roles of Tape 1 and Tape 2 from the next cycle.
- The new memory to be stored on Tape 0 is mem_{k+1} and the conjecture is e_k .

It is easy to see that in each cycle, the time spent is proportional to $|mem_k| + |w_k| + |t|$ and thus linear in the length of the longest word seen so far (plus a constant); note that mem', e', e_k are also bounded by that length (plus a constant). Furthermore, in the simulation of M , each input word to N is given to M infinitely often. Hence N learns each language from the class \mathcal{L} . \square

Open Problem 3.5. It is unknown whether one can learn every in principal learnable automatic class using an automatic learner augmented by only one work tape.

Further investigations deal with the question what happens if one does not add further work tapes to the learner but uses other methods to store memory. Indeed, the organisation in a tape is a bit awkward and using a queue solves some problems. A queue is a tape where one reads at one end and writes at the opposite end, both the reading and writing heads are unidirectional and cannot overtake each other. Tape 0 satisfies the same constraints as

in the model of additional work tapes and one also has the constraint that in each cycle only linearly many symbols (measured in the length of the longest datum seen so far) are stored in the queue and retrieved from it.

Theorem 3.6. *Every learnable automatic family \mathcal{L} has a linear-time learner using one additional queue as a data structure.*

Proof. The learner simulates an automatic learner M for \mathcal{L} using fat text, in a way similar to that done in Theorem 3.4. Let M in the k -th step map (mem_k, w_k) to (mem_{k+1}, e_k) for M 's memory mem_k .

For ease of presentation, the contents of Tape 0 is considered as consisting of a convolution of 4 items (rather than 2 items, as considered in other parts of the paper). At the beginning of a cycle the linear-time learner N has $\text{conv}(v_k, -, mem_k, -)$ on Tape 0 where v_k is the current datum, mem_k the archived memory of M and “-” refers to irrelevant or empty content. In the k -th cycle, the linear-time learner N scans four times over Tape 0 from beginning to the end and each time afterwards returns to the beginning of the tape:

- (1) Copy v_k from Tape 0 to the write-end of the queue;
- (2) Read a word from the read-end of the queue, call it w_k , and update Tape 0 to $\text{conv}(v_k, w_k, mem_k, -)$;
- (3) Copy w_k from Tape 0 to the write-end of the queue;
- (4) Simulate M on Tape 0 in order to map (mem_k, w_k) to (mem_{k+1}, e_k) and update Tape 0 to $\text{conv}(v_k, w_k, mem_{k+1}, e_k)$.

It can easily be verified that this algorithm permits to simulate M using the data type of a queue and that each cycle takes only time linear in the length of the longest datum seen so far. Thus, N learns \mathcal{L} . \square

A further data structure investigated is the provision of additional stacks. Tape 0 remains a tape in this model and has still to obey to the resource-bound of not being longer than the longest word seen so far (plus a constant). Theorems 3.3 and 3.4 work also with one and two stacks, respectively, as the additional work tapes are actually used like stacks.

Theorem 3.7. *There is an automatic class which can be learnt with one additional stack but not by an automatic learner. Furthermore, every learnable automatic class can be learnt by a learner using two additional stacks.*

Furthermore, the next result shows that in general one stack is not enough; so one additional stack gives only intermediate learning power while two or more additional stacks give the full learning power. The class witnessing the separation contains only finite sets.

For information on Kolmogorov complexity, the reader is referred to standard text books [5, 9, 24, 25]. The next paragraphs provide a brief description of the basic concepts.

Consider a Turing machine U which computes a partial-recursive function from $\{0, 1\}^* \times \{0, 1\}^*$ to $\{0, 1\}^*$. The first input to U is also referred to as a program. Machine U is universal iff for every further machine V , there is a constant c such that, for every (p, y) in the domain of V , there is a q , which is at most c symbols longer than p , satisfying $U(q, y) = V(p, y)$. Fix a universal machine U . Now the conditional Kolmogorov complexity $C(x|y)$ is the length of the shortest program p with $U(p, y) = x$; the plain Kolmogorov complexity $C(x)$ is $C(x|\varepsilon)$. Note that, due to the universality of U , the values of $C(\cdot)$ can only be improved by a constant (independent of x, y) when changing from one universal machine to another one. In some cases below, $C(x|y_1, y_2, \dots, y_r)$ is the conditional Kolmogorov

complexity when given an r -tuple (y_1, y_2, \dots, y_r) where r might vary; such a tuple can be coded up in any way which permits to identify the parts uniquely, as automaticity is not required, the coding $0^{|y_1|}1y_10^{|y_2|}1y_2 \dots 0^{|y_r|}1y_r$ would do it.

If f is a partial-recursive function then there is a constant c with $C(f(x)|y) \leq C(x|y) + c$ for all x in the domain of f . In particular, if one can find a way to describe the strings x in a set A by binary strings p_x such that some algorithm can compute each $x \in A$ from the corresponding description p_x , then $C(x) \leq |p_x| + c$ for some constant c and all $x \in A$. For this reason, one often says that x can be described by n bits when the corresponding p_x can be chosen to have n bits.

Theorem 3.8. *The class of all $L_e = \{x \in \{0, 1\}^{|e|} : x \neq e\}$ with $e \in \{0, 1\}^* \cup \{2\}^*$ cannot be learnt by a linear-time learner using one additional stack.*

Proof. Assume that the linear-time learner M using one stack, in addition to the base tape, is given. In order to find languages not learnt by M , one focuses on L_e where the parameter $n = |e|$ is large; in addition one considers only n of the form $k + 2^k$ for some k ; this parameter k and $m = 2^k$ will play some role in the arguments below. Note that all data-items in L_e have the length n . For $i \in \{1, 2, \dots, m\}$, let x_i be a string of length n such that the Kolmogorov complexity $C(x_1x_2 \dots x_m)$ is at least $(n - k)m$ and the first k bits of each x_i is the binary bit representation of $i - 1$. Furthermore, assume that c is a constant so large that the Kolmogorov complexity of the content of Tape 0 (which can be assumed to be always n symbols long, since all data have length n , but which can use more than two alphabet symbols) is at most cn and that the stack can, in each round, pull or push up to cn symbols, where the stack alphabet has at most 2 symbols (one can code up a larger alphabet in binary and choose the constant c sufficiently large to absorb the extra amount of storage). Hence, in each cycle, what the machine does depends on the content of Tape 0 (worth cn bits) and on the top cn symbols of the stack (worth cn bits). Furthermore, assume that all words of length n different from x_1, x_2, \dots, x_m have already been presented to the learner and let α denote the content of Tape 0 and $u\beta$ denote the content of the stack where β are the top cn symbols (or less if u is the empty word). Below one considers the behaviour/configuration of the learner when it is presented with further inputs and one considers $(\alpha, u\beta)$ as the initial configuration of the learner for this purpose. Below, the configuration of the learner, at any stage before reading the next input, is denoted by (\cdot, \cdot) , where the first argument is the content of the tape and the second argument is the content of the stack.

Intuitively, as the x_j 's are complex, the learner needs to store them on the stack when it receives them (otherwise, it would lose information about which x_j 's it has seen). This forces the stack to grow larger and larger and prevents the learner from accessing earlier stored data on the stack, thus making the earlier stored information useless. This allows to show that some language L_e is not learnable by M .

Claim 3.11 gives a permutation $x_{i_1}, x_{i_2}, \dots, x_{i_m}$ of x_1, x_2, \dots, x_m such that M does not touch any, but the top $6(c + 1)^2n$ symbols of u , on input $x_{i_1}, x_{i_2}, \dots, x_{i_m}$. Claim 3.12 uses this claim to show that on some sequence σ of x_i 's, M reaches a configuration $(\alpha', vw'\beta')$, with $|\beta'| = cn$ and v being same as u except for the top $6(c + 1)^2n$ symbols removed, where the learner never touches v on the input σ , and for any future input involving x_i 's never touches vw' . This, then allows to claim in Claim 3.13 that the learner cannot learn some L_e . Claims 3.9 and 3.10 are used in proving the above claims. Now the five claims about the configuration of M are proven formally.

Claim 3.9. There do not exist two distinct input sequences of words of length n , one containing an x_i and one not containing x_i , ending up in the same configuration $(\alpha', v\beta')$ where β' has at most length cn and v has not been touched (that is, starting from configuration $(\alpha, u\beta)$, the initial portion v of $u\beta$ above was never at the top of the stack during the processing of any of the two sequences).

Assume by way of contradiction that this claim fails, that is, there are two such sequences σ and σ' . Then one can bring the learner into the configuration $(\alpha', v\beta')$ by either of the sequences and thereafter feed the learner with the x_j with $j \neq i$, and then with a string of length n , different from x_i , forever. The convergence behaviour of the learner, in both cases, is the same as the configuration $(\alpha', v\beta')$ is independent of the sequence σ or σ' by which the learner reached it; from then onwards the learner receives, in both cases, the same data and conjectures the same hypotheses, as in both cases they are based on the same data, Tape 0 and stack. In one case the learner has to learn $L_{x_i} = \{0, 1\}^n - \{x_i\}$ and in the other case the learner has to learn $L_{2^n} = \{0, 1\}^n$; thus the learner can learn at most one of these two sets. This completes the proof of the claim.

Claim 3.10. There is no input sequence $(x_{i_1}, x_{i_2}, \dots, x_{i_\ell})$ and no splitting of u into vw such that M , after reading these inputs, is in a configuration of the form $(\alpha', v\beta')$ with $|\beta'| \leq cn$ and without having pulled and pushed back any symbols of v and with the conditional Kolmogorov complexity satisfying $C(x_{i_1}x_{i_2} \dots x_{i_\ell} | \alpha, w\beta, i_1, i_2, \dots, i_\ell) \geq (c+1)^2n$.

For a proof of the claim, assume by way of contradiction that there is such an input sequence $(x_{i_1}, x_{i_2}, \dots, x_{i_\ell})$. Then there is a partial-recursive function f such that f , given $(\alpha, w\beta, i_1, i_2, \dots, i_\ell, \alpha', \beta')$, finds a sequence $y_{i_1}, y_{i_2}, \dots, y_{i_\ell}$ such that $y_{i_j} = y_{i_{j'}}$ whenever $i_j = i_{j'}$, $y_{i_j} \in \{0, 1\}^n$ for all j , y_{i_j} having the first k bits being the binary representation of i_j and M pulling on these inputs the symbols belonging to $w\beta$ without touching those of v and ending up in the configuration $(\alpha', v\beta')$. Note that one does not need to know v for this search, hence the search depends only on the inputs given to f and returns an input sequence such that its Kolmogorov complexity given $(\alpha, w\beta, i_1, i_2, \dots, i_\ell)$ is at most that of (α', β') , that is, below $(c+1)^2n$ (assuming that n is sufficiently large). It follows that at least one y_{i_j} differs from x_{i_j} ; furthermore, no other $y_{i_{j'}}$ can be equal to x_{i_j} by the rules that each $y_{i_{j'}}$ encodes $i_{j'}$ in the first k bits and equals to y_{i_j} whenever $i_{j'} = i_j$. However, this would contradict Claim 3.9. This completes the proof of Claim 3.10.

Claim 3.11. There is a permutation $(x_{i_1}, x_{i_2}, \dots, x_{i_m})$ of (x_1, x_2, \dots, x_m) such that the splitting $vw = u$ with either $|w| = 6(c+1)^3n$ or $|w| < 6(c+1)^3n \wedge |v| = 0$ satisfies that M on input $(x_{i_1}, x_{i_2}, \dots, x_{i_m})$ never touches the symbols in v .

Let vw be the given splitting of u . If $|w| < 6(c+1)^3n$ then v is empty and nothing needs to be proven; thus assume that $|w| = 6(c+1)^3n$.

Now define $y_k = x_1x_2 \dots x_m$ and inductively for $\ell = k-1, k-2, \dots, 1$, split $y_{\ell+1}$ at the middle into two equal parts y_ℓ and z_ℓ (both of length $2^\ell n$) such that $C(y_\ell | k, \alpha, w\beta) \geq C(z_\ell | k, \alpha, w\beta)$. Note that there is a unique permutation of the form $(x_{i_1}, x_{i_2}, \dots, x_{i_m})$ of (x_1, x_2, \dots, x_m) such that

$$x_{i_1}x_{i_2} \dots x_{i_m} = y_1z_1z_2 \dots z_{k-1}.$$

Note that i_2, i_3, \dots, i_m can be computed from i_1 . Note that $C(y_k | k, \alpha, w\beta) \geq (n-2k)m$ (for k and $n = k+2^k$, $m = 2^k$ being sufficiently large) for the following reasons: $C(y_k) \geq (n-k)m$; $C(y_k | k, \alpha, w\beta) \geq C(y_k) - C((k, \alpha, w\beta)) - k$; $C(k, \alpha, w\beta) + k \leq 8(c+1)^3n \leq km/2$.

By induction one can see that $C(y_\ell | k, \alpha, w\beta) \geq (n-2k)m \cdot 2^{\ell-k} - k$ for all ℓ whenever

k, n are sufficiently large; note that y_ℓ is the more complex half of $y_{\ell+1}$ and therefore has by induction hypothesis at least the complexity $(n - 2k) \cdot m \cdot 2^{\ell+1-k}/2 - k/2$ minus some constant which can be brought into the form $(n - 2k)m \cdot 2^{\ell-k} - k$ by assuming that $k/2$ is larger than the corresponding constant.

Furthermore, the values i_1, i_2, \dots, i_m can be computed from k and i_1 , hence one can represent i_1, i_2, \dots, i_h by i_1 and h and k . Hence

$$C(y_\ell | \alpha, w\beta, i_1, \dots, i_h) \geq (n - 2k)m \cdot 2^{\ell-k} - 5k$$

for any h with $2^\ell \leq h < 2^{\ell+1}$. There are two cases for each h with $2^\ell \leq h < 2^{\ell+1}$:

First, $3(c+1)^2 n > |y_\ell|$. Then $h < 6(c+1)^2$ and, on input $(x_{i_1}, x_{i_2}, \dots, x_{i_h})$, the learner can have pulled at most $6c(c+1)^2 n$ symbols from the stack; hence it has neither touched v nor the bottom cn symbols of w .

Second, $3(c+1)^2 n \leq |y_\ell|$. Then

$$\begin{aligned} C(y_\ell | \alpha, w\beta, k, i_1, i_2, \dots, i_h) &\geq (n - 2k)m \cdot 2^{\ell-k} - 5k \geq \\ 3(c+1)^2(n - 2k) - 5k &> 2(c+1)^2 n. \end{aligned}$$

Assuming that k and $n = 2^k + k$ are sufficiently large, one obtains

$$C(x_{i_1} x_{i_2} \dots x_{i_h} | \alpha, w\beta, i_1, i_2, \dots, i_h) \geq (c+1)^2 n.$$

Thus, using Claim 3.10 it follows that for all $h \in \{6(c+1)^2, 6(c+1)^2 + 1, \dots, m\}$ there are at least cn symbols in the stack above v after reading $x_{i_1}, x_{i_2}, \dots, x_{i_h}$.

Hence, using above cases, one can conclude by induction on h that the symbols in v are not touched while processing the input $(x_{i_1}, x_{i_2}, \dots, x_{i_m})$.

Claim 3.12. Split u into vw as in Claim 3.11. There is a sequence of all x_i , perhaps with repetitions, such that after reading this sequence M is in a configuration $(\alpha', vw'\beta')$, with $|\beta'| = cn$, such that for all further inputs from x_1, x_2, \dots, x_m , M does not touch the symbols on the part of the stack denoted by vw' .

Assuming that this sequence does not exist, one could use the sequence given in Claim 3.11 to remain above v in the stack until all symbols are passed and then one could feed some sequence of x_i until all but at most cn symbols above v are used up; that is, one would be in a configuration of the form $(\alpha', v\beta')$ with $|\alpha'| = n$ and $|\beta'| \leq cn$. Hence one can, given $(\alpha, w\beta)$ and (α', β') search a tuple (y_1, y_2, \dots, y_m) such that each y_i starts with a binary number of length k representing $i - 1$ and each y_i has n bits and there is a sequence of inputs drawn from this tuple on which the configuration of M with $(\alpha, vw\beta)$ changes to $(\alpha', v\beta')$ without touching v . The first tuple (y_1, y_2, \dots, y_m) of this type found by searching has Kolmogorov complexity at most $8(c+1)^3 n$ (obtained by coding the inputs $k, \alpha, w\beta, \alpha', \beta'$ and the routine for the search programme) which is less than $(n - k)m$, the lower bound on the Kolmogorov complexity of $x_1 x_2 \dots x_m$, for sufficiently large k, m, n . Therefore some y_i differs from x_i and therefore one can reach the configuration $(\alpha', v\beta')$ from $(\alpha, vw\beta)$ by either having seen x_i or not having seen x_i . It follows from Claim 3.9 that this cannot occur, hence there is some minimal extension $w'\beta'$ of v such that $|\beta'| = cn$ and when reading any sequence of the data x_1, x_2, \dots, x_m after having reached the configuration $(\alpha', vw'\beta')$, it will not touch vw' in the stack, that is, all future activity depends only on α' and β' .

Claim 3.13. M fails to learn some language of the form $\{0, 1\}^n - \{x_i\}$ or $\{0, 1\}^n$.

Let $u, v, w, \alpha, \beta, \alpha', \beta', w'$ as in Claim 3.12. One can now show that there is a tuple (y_1, y_2, \dots, y_m) with $|y_i| = n$ and y_i extending the k -bit representation of $i - 1$ such that

M when fed with some input-sequence taken from $\{y_1, y_2, \dots, y_m\}$ ends up in a configuration of the form $(\alpha', vw''\beta')$ without touching v and this configuration is computed from $(\alpha, w\beta, \alpha', \beta')$; as in Claim 3.12 one can argue that some $y_i \neq x_i$. Now one can feed all the $x_j \neq x_i$ into M for the configurations $(\alpha', vw'\beta')$ and $(\alpha', vw''\beta')$, respectively, for both in the same way and in a loop repeated forever. In both cases the learner M either converges to the same index or does not converge, but in one case the text which M has received is a text for $\{0, 1\}^n$ and in the other case it is a text for $\{0, 1\}^n - \{x_i\}$. Hence M fails to learn at least one of these two sets. \square

4. RELAXING THE TIMING CONSTRAINTS

In this section, it is investigated how the learning power improves if the severe restrictions on work Tape 0 or the computation time are a bit relaxed. The next result shows that, if one allows a bit more than just linear time, then one can learn, using one work tape, all learnable automatic classes of infinite languages. The result could even be transferred to families of arbitrary r.e. sets as the simulated learner is an arbitrary recursive learner. Intuitively, think of f in the following theorem as a slowly growing function.

Theorem 4.1. *Assume that $\{L_e : e \in I\}$ is an automatic family where every L_e is infinite and M is a recursive learner which learns this family. Furthermore, assume that f, g are recursive functions with the property that $f(n) \geq m$ whenever $n \geq g(m)$ (so g is some type of inverse of f). Then there is a learner N which learns the above family, using only one additional work tape, and satisfies the following constraint: if n is the length of the longest example seen so far, then only the cells number $1, 2, \dots, n$ of Tape 0 can be non-empty and the update time of N in the current cycle is $O(n \cdot f(n))$.*

Proof. The main idea of the proof is that one constructs a learner which splits Tape 1 into four tracks for archivation; the learner usually uses Track 1; in irregular intervals, the learner returns from its current position to the origin of Tape 1 and uses Track 2 for archiving the examples which come up during this “return to the origin” until it reaches the old data on Tracks 2 and 3. When this happens, the old data found there consist only of words up to length m (where m is sufficiently small compared to the current word length n) and the learner can compress the data in Tracks 1, 2 and 3 into a list α (to be maintained on Tape 0); α will contain, for each word w up to length m occurring in the input, at most one copy (which gives a corresponding length bound on the length of α). Once the compression is completed, the learner returns to the forward mode using the one left over free track for this purpose. The key idea is to “space out” the visits to the origin such that, for m being the length of the longest datum seen up to the end of the last visit, m is so much smaller than the current n that $2^{m+1} \cdot (m+1) \leq f(n)$; this allows all the data which was archived up to the end of the previous visit to be compressed into a string of length up to $f(n)$ and the update of this compressed memory can, in each round, be done in time $O(f(n) \cdot n)$.

The description below gives a more precise description of the update protocol. As the memory has only to be bounded by the length of the longest datum seen so far plus some constant, one can assume without loss of generality that n is at least 1.

On Tape 0, as memory, the learner N archives the convolution of variables $\alpha, \beta, \gamma, e, 0^m, 0^n$ with the following meaning.

- 0^n represents in unary the length of the longest word seen so far and 0^m is an old value of 0^n ; initially m and n are 1 (not 0).

- The variable α is, during the runtime, only modified by appending symbols at the end and will in the limit consist of a one-one text of all the words occurring in the language to be learnt; the words on α are separated by a special character. For example, $\alpha = \#00\#\#0101\#11111\#$ would represent a beginning of a text consisting of 00, ε , 0101 and 11111. Furthermore, each of the words in α would be of length at most m .
- The variable β is the current configuration of a computation to determine $g(2^{m+1} \cdot (m+1))$ (in unary); this configuration is updated whenever the length and time constraints permit and the next configuration is shorter than 0^n , until the computation finishes.
- The variable γ is a configuration of M , while processing the initial part α of a text for the input language; note that this configuration includes the memory of M and the portion of α it has read. In each cycle this configuration is updated by one more step of the computation, unless the input α is currently exhausted (that is, M would like to read a symbol which is not yet there) or the length of the configuration becomes longer than 0^n .
- The variable e is the last completed conjecture of M and updated whenever the configuration γ of M contains a new value to be output.

In each cycle, the learner N would archive the current input x on the work tape at a position near to the current one (that is, the input position has to be reached in linear time) and N would furthermore update the values of $\beta, \gamma, e, 0^m, 0^n$ on Tape 0 (α is updated only during some cycles, see below).

In order to be able to save all required information on the work Tape 1, the tape content is modeled as having four tracks. Usually, only Track 1 is used for appending new information at the end of the tape and Track 4 is used for making sure that computations of the variables of Tape 0 meet the time-bound. Tracks 2 and 3 are used to store data during cycles when some special operations are needed to transfer data from Tape 1 to the memory α in Tape 0. Furthermore, initially $m = 1$.

When β shows that the computation of $g(2^{m+1} \cdot (m+1))$ has terminated, and the observed examples are so long that $n \geq g(2^{m+1} \cdot (m+1))$ then the learner enters the phase to do special operations (for next several cycles, as many as needed). Note that eventually this happens for every value of m , as the input language is infinite (assuming it is from \mathcal{L}). In each cycle during this special phase, from its current position at the end of Tape 1 back to the origin \boxplus , N will transfer/copy all stored words in Tape 1 of length at most m , which are not already in α , to α . During this process, the older words stored in Tracks 2 and 3 may be erased (but not lost, as they have already been copied to α , as each of them are of length at most m). The new input words received during this phase are copied in Tracks 2 and 3 (see below). Note that a concatenation of all words up to length m is at most $2^{m+1} \cdot (m+1)$ long (including separating symbols) and hence $|\alpha| \leq 2^{m+1} \cdot (m+1)$ whenever α consists only of copies of words up to length m appearing in the language to be learnt and each such word appears at most once in α .

Now, it is described how special operations are done in the special phases, see also Figure 2 for a rough summary of the handling of old and new data in each cycle. When going back on Tape 1, N will do the following for all words w archived in Tracks 1, 2, 3 starting from the current position up to $|x| + 1$ positions left of the current position (here one also considers w that might only partially overlap with the cells in positions between the current position and $|x| + 1$ to the left of the current position; recall that x is the current input data to the learner): if $|w| \leq m$ then w is compared with all words in α and in the case that it does not coincide with any archived word in α , $w\#$ is appended at the end of α ; note that all words archived in the Tracks 2 and 3 have at most the length m . For each word

Mode	Usual	Backward Special	Forward Special
Old Data Before	—	In Tracks 1, 2 and 3	In Tracks 1 and 2
Old Data After	—	Into Base Tape and Track 1	Remains unchanged
New Data	Into Track 1	Into Track 2	Into Track 3

Figure 2: Handling of data at head position of Tape 1. In backward special mode, old short data is recorded into the base tape and old long data remains in Track 1.

w , this operation needs time $O(|\alpha| \cdot |w|)$. Note that w has at most length m and α at most length $2^{m+1} \cdot (m+1)$, giving an overall bound of $O(2^{m+1} \cdot (m+1) \cdot |w|)$ for the processing of each word w . Furthermore, the concatenation of all these words archived one after another has length at most $3|x| + 3m$; so one can conclude that the whole operation needs time $O(2^{m+1} \cdot (m+1) \cdot n)$ which is $O(f(n) \cdot n)$ as $g(2^{m+1} \cdot (m+1)) \leq n$. Furthermore, all w in Tracks 2 and 3 overlapping with the space between the current position in Tape 1 and the cell at position $|x|$ left of the current position before the start of the cycle are cleared away as these w all have at most the length m . After the clearance, x will be archived in Track 2 (where a special symbol outside the alphabet used for the archivation data is used to fill up blank spaces, if needed) and the current position moves by $|x| + 1$ to the left. This is done until the origin \boxplus is reached. At this point, Track 3 is empty and can be used to archive the incoming data in a similar way while the Turing machine moves back from \boxplus to end of used part of Tape 1. When returning to the usual archivation mode, m is updated to be the current value of n so that all words archived in Tracks 2 and 3 are again having at most length m . From then onwards, one waits until so much data has been observed such that the computation of $g(2^{m+1} \cdot (m+1))$ has terminated and gives a value below (the new value of) n .

One can see from this description that, when learning an infinite language, eventually all words observed will be appended to α and M will be simulated on the resulting one-one text of the language to be learnt. Thus, M will eventually stabilise on some index e , which will be taken over as output when the corresponding computation has terminated and n is larger than $|e|$. This shows that N follows the simulated learner M and therefore N learns the class to be learnt. \square

Pitt's original result [28] on linear time learners did not measure the time in the size of the largest example seen so far, but in the size of the overall amount of examples seen so far. So the next two results deal with the question of the additional learning power provided by one work tape or one stack when the learner can use a Tape 0 of length n and run in time linear in n where n is logarithm of the number of data seen so far plus the length of the longest example seen so far; hence n increases, though slowly, when a datum is presented multiply.

Note that in the proof of Theorem 4.1, the main reason to use infinite languages and strings of larger and larger length n , was to be able to transfer all stored data of length m onto α . This can also be done if instead of the length n , the unbounded growing number of examples seen so far is used as a parameter to allow the time needed to do the transfer (in which case additionally, one can make α a fat text). For this, one needs to keep track of some earlier maximal length m' and number of items n' (including counting the multiple copies, in case they are there) so that $2 \cdot m' \cdot n'$ bounds the overall length of all examples

stored in Tracks 2 and 3. When the number of examples seen so far, n , is larger than the current length of α plus $2 \cdot m' \cdot n'$, one can then start going back, copying new data in Track 2 until one reaches the point where the earlier data in Tracks 2 and 3 are stored. At this point one moves all data in Tracks 2 and 3 to the end of α which is stored in Tape 0 and then starts moving forward on Tape 1 again, copying new data into Track 3 until one reaches the end of recorded part of all the tracks. At this point one can consider Track 1 and Track 2 as old recorded data (earlier roles played by Tracks 2 and 3) and continue recording data in Track 3 up to the point when the learner has seen enough examples so as to copy the data in Tracks 1 and 2 to Tape 0. Continuing in this way, one can copy all data to α in Tape 0 eventually and use the data in Tape 0 to simulate an automatic learner on fat text by cycling through the examples archived in α .

This allows to show the following result; its proof is similar to Theorem 4.1 and the details are omitted.

Theorem 4.2. *Let n be the logarithm of the number of data seen so far plus the length of the longest example seen so far and consider a learner which can store in Tape 0 information of length n and can access one additional work tape, with update time in each cycle being linear in the corresponding n . Then such a learner can learn every learnable automatic family.*

The previous and the next result compute the parameter n of the update time and length of Tape 0 in the same way. While the previous result showed that one additional work tape is sufficient for full learning power under the corresponding linear time model, the next result shows that one additional stack is insufficient for full learning power.

Theorem 4.3. *Let n be the logarithm of the number of data seen so far plus the length of the longest example seen so far and consider a learner which can store in Tape 0 information of length n and can access one additional stack, with update time in each cycle being linear in the corresponding n . Then such a learner fails to learn the class \mathcal{L} of all set $L_e = \{0, 1\}^* - \{e\}$ where the indices e range over $\{0, 1\}^*$.*

Proof. Assume by way of contradiction that such a learner M for \mathcal{L} exists.

Intuitively, the idea of the proof is that if the learner gets complex strings (relative to the position), then it has to store it in the stack. Thus, if it gets complex strings in odd positions of the text, and even positions of the text are filled with simple strings (to form a complete text for some target language), then the learner has to push (codings of) the complex strings on the stack and is not able to look at these pushed symbols in later computation. This allows to construct two such texts for different languages in the class on which eventually the learner behaves in the same way (see Claim 4.8, and then the arguments after this claim); thus the learner can learn at most one of these two sets. Claims 4.4 to 4.7 are combinatorial claims based on Kolmogorov complexity, needed for proving Claim 4.8. Now the formal proof is given.

Let $\text{bin}(m)$ denote the binary representation of m using $\log(m+2)$ bits where, for $k \geq 1$, $\log(k)$ is the downrounded logarithm of base 2, that is, the maximal integer h with $2^h \leq k$.

Claim 4.4. Suppose σ and τ are two finite sequences over $\{0, 1\}^*$ such that $\text{range}(\sigma) - \text{range}(\tau) \neq \emptyset$, $\text{range}(\tau) - \text{range}(\sigma) \neq \emptyset$, and M has the same Tape 0 content and stack content after processing either σ or τ . Then, M does not learn \mathcal{L} .

To show the above claim, let $w \in \text{range}(\sigma) - \text{range}(\tau)$ and $w' \in \text{range}(\tau) - \text{range}(\sigma)$. Let T' be a text for $\{0, 1\}^* - \{w, w'\}$. Then, M has the same convergence behaviour (that is it

either diverges or converges to the same conjecture) on texts $\sigma T'$ and $\tau T'$, which are texts for $L_{w'}$ and L_w respectively. Thus, M fails to learn at least one of these languages and thus fails to learn \mathcal{L} . This completes the proof of the claim.

For any recursive text T of any language satisfying $|T(m)| \leq \log(m+2)$ for all m , define R_T (using an oracle for the halting problem K) as follows: Let $R_T(2m) = T(m)$ and $R_T(2m+1)$ be the string x of length $16(\log(m+2))$ ending with $\text{bin}(m)10^{\log(m+2)-1}$ which maximises $C(x|R_T(0)\#R_T(1)\#\dots\#R_T(2m))$.

Claim 4.5. Let T be a recursive text satisfying $|T(m)| \leq \log(m+2)$ for all m . Then the following statements hold:

- (a) $R_T(2j+1)$ are pairwise distinct for different j ;
- (b) For each m , $R_T(2m+1) \notin \{T(i) : i \leq m^2\}$;
- (c) For each m , $C(R_T(2m+1)|R_T(0)\#R_T(1)\#\dots\#R_T(2m)) \geq 14 \log(m+2)$.

Part (a) follows by definition. Part (b), follows by definition of $R_T(2m+1)$ and the fact that $16 \log(m+2) > \log(m^2+2)$. For part (c) note that there exists a string x of length $16 \log(m+2)$, which ends in $\text{bin}(m)10^{\log(m+2)-1}$, with Kolmogorov complexity (given $R_T(0)\#R_T(1)\#\dots\#R_T(2m)$) at least $14 \log(m+2)$. As $R_T(2m+1)$ is most complex such string x , part (c) follows.

Claim 4.6. There exists a constant c_2 such that the following holds for $m \geq c_2$. Suppose T is a recursive text satisfying $|T(i)| \leq \log(i+2)$ for all i . Then, $C(R_T(2m+1)\#R_T(2m+2)\#\dots\#R_T(2m+2k-1)|R_T(0)\#R_T(1)\#\dots\#R_T(2m)) \geq k \cdot \log(m+2)$.

To see that the claim holds, note that for some constant c_1 , for all $x, y \in \{0,1\}^*$, $\sigma \in \{0,1,\#\}^*$, $C((x,y)|\sigma) \geq C(x|\sigma) + C(y|\sigma\#x) - c_1$, see [24]. Thus, for all large enough m ,

$$\begin{aligned} & C(R_T(2m+1)\#R_T(2m+2)\#\dots\#R_T(2m+2k)|R_T(0)\#R_T(1)\#\dots\#R_T(2m)) \\ & \geq \sum_{i=m}^{i=m+k-1} [14 \log(i+2) - c_1] \geq k \log(m+2) \end{aligned}$$

(where the second last inequality follows from Claim 4.5(c)).

Let U_i^T and V_i^T denote the Tape 0 content and stack content of M after processing $R_T(0), R_T(1), \dots, R_T(2i)$.

Claim 4.7. There exists a constant c_3 such that, for m and k greater than c_3 , with $m+k \leq m^2$, the following holds. Suppose T is a recursive text satisfying $|T(i)| \leq \log(i+2)$ for all i . Furthermore suppose that T is computed by a program of Kolmogorov complexity less than $6 \log(m+2)$. Then,

$$C((U_{m+k}^T, V_{m+k}^T)|(U_m^T, V_m^T, k)) \geq \frac{k \log(m+2)}{3}.$$

To show that the claim holds, suppose m is large enough as required for Claim 4.6. Suppose $C(U_{m+k}^T, V_{m+k}^T|U_m^T, V_m^T, k) < \frac{k \log(m+2)}{3}$. Note that by Claim 4.5 (b), for i with $m \leq i < m+k$, $R_T(2i+1)$ does not belong to $T(0), T(1), \dots, T(m+k)$. Now, given a program for T and $U_m^T, V_m^T, k, U_{m+k}^T, V_{m+k}^T$, one can construct $w_{2m+1}, w_{2m+3}, \dots, w_{2m+2k-1}$ such that, for i with $m \leq i < m+k$,

- (i) w_{2i+1} ends in $\text{bin}(i)10^{\log(i+2)}$,
- (ii) $w_{2i+1} \notin \{T(0), T(1) \dots T(m+k)\}$,

- (iii) M starting with Tape 0 content U_m^T and stack content V_m^T , on input sequence $w_{2m+1}T(m+2)w_{2m+3}\dots w_{2m+2k-1}T(m+k)$, ends in Tape 0 content being U_{m+k}^T and stack content V_{m+k}^T .

Note that U_m^T and V_m^T can be computed using $R_T(0)\#R_T(1)\#\dots\#R_T(2m)$. Thus, the expression

$$C(w_{2m+1}\#R_T(2m+2)\#w_{2m+3}\#R_T(2m+4)\dots\#w_{2m+2k-1}|R_T(0)\#R_T(1)\#\dots\#R_T(2m))$$

is bounded by $C(U_{m+k}^T, V_{m+k}^T|U_m^T, V_m^T, k) + 2\log(k+2) + 12\log(m+2) + c' \leq \frac{k\log(m+2)}{3} + 2\log(k+2) + 12\log(m+2) + c'$ for some constant c' . However, by Claim 4.6,

$$C((R_T(2m+1)\#R_T(2m+2)\#\dots\#R_T(2m+2k-1))|R_T(0)\#R_T(1)\#\dots\#R_T(2m)) \geq k \cdot \log(m+2).$$

Thus, for large enough m , there are some h, w, w' satisfying $m \leq h < m+k$, $w' = w_{2h+1} \neq w = R_T(2h+1)$ and $w, w' \notin \{T(0), T(1), \dots, T(m+k)\}$. But, then by Claim 4.4 and Claim 4.5(a), M does not learn \mathcal{L} . Hence, $C((U_{m+k}^T, V_{m+k}^T)|(U_k^T, V_k^T, k)) \geq \frac{k \cdot \log(m+2)}{3}$. This proves Claim 4.7.

Claim 4.8. There exists a constant c_5 such that for large enough m and $m+k \leq m^2$ the following holds. Suppose T is a recursive text satisfying $|T(i)| \leq \log(i+2)$ for all i , and T is computed by a program of Kolmogorov complexity less than $6\log(m+2)$. Then, while processing $R_T(2m+1)R_T(2m+2)\dots R_T(2m+2k)$,

- (a) the part of stack consisting of V_m^T , except for the top $c_5 \log(m+2)$ symbols, is never removed and
 (b) $|V_{m+k}^T| \geq (|V_m^T| - c_5 \log(m+2)) + \frac{k \log(m+2)}{4}$.

To show the claim, consider $k \leq m^2 - m$. Now, $|U_{k+m}^T| = O(\log(k+m+2))$, and thus $C(U_{k+m}^T|U_m^T, V_m^T, k) = O(\log(k+m+2))$. Furthermore, if $V_{k+m}^T = vw$, for some longest prefix v of V_m^T , then the length of the deleted portion of V_m^T (that is $|V_k^T| - |v|$), can be at most $k \cdot \log(k+m+2) \cdot c'$, for a constant c' ; this can be coded using $\log(k+2) + \log \log(k+m+2) + c'$ bits. Hence, by Claim 4.7, for some constant c_4 , $|w| \geq \frac{k \cdot \log(m+2)}{3} - c_4[\log(m+2) + \log(k+2)] \geq \frac{k \cdot \log(m+2)}{4}$, for $m, k \geq c''$, for some constant c'' . Thus, for large enough k , the length of w above is larger than what can be removed from the stack in one cycle. It follows that, for some constant c_5 , the machine M , on input $R_T(2m+1), R_T(2m+2), \dots, R_T(2m+k)$, does not remove symbols from V_m^T , except maybe for up to $c_5 \log(m+2)$ symbols from the top. This proves part (a). Part (b) follows, by using the length of w above. This completes the proof of Claim 4.8.

Using (a) and (b) of the above claim, it follows that for large enough m , for each recursive text T of some subset of $\{0, 1\}^*$ and T having a program shorter than $6\log(m+2)$, M on $R_T(0), R_T(1), \dots, R_T(2m)$ will, when processing subsequent data from R_T , never remove symbols from V_m^T except maybe for the top $c_5 \cdot \log(m+2)$ symbols.

Now, given a program for T , m , $R_T(0)\#R_T(1)\#\dots\#R_T(2m)$, U_{2m}^T and the topmost $c_5 \cdot \log(2m+2)$ symbols of V_{2m}^T , one can compute a sequence σ of length $4m+1$ such that,

- (i) $\sigma(s) = R_T(s)$, for $s \leq 2m$,
 (ii) for all $i \leq 2m$, $\sigma(2i) = T(i)$ and
 (iii) for all $i < 2m$, $\sigma(2i+1)$ ends with $\text{bin}(i)10^{\log(i+2)}$ and is of length $16\log(i+2)$.

(iv) M after processing σ has Tape 0 content U_{2m}^T and the top $c_5 \log(2m+2)$ symbols of the stack are same as the top $c_5 \log(2m+2)$ symbols of V_{2m}^T .

Now, $C(\sigma|R_T(0)\#R_T(1)\#\dots\#R_T(2m))$, is at most $c_6 \cdot \log(m+2)$, for some constant c_6 , as it was constructed from a description of m and a description of the top $c_5 \log(2m+2)$ stack symbols and $O(\log(2m+2))$ symbols of U_{2m}^T . On the other hand, $C(R_T(2m+1)\#R_T(2m+2)\#\dots\#R_T(2m+2m)|R_T(0)\#R_T(1)\#\dots\#R_T(2m)) \geq m \cdot \log(m+2)$ by Claim 4.6.

Hence, fix a text $T = T_0$ of the nonempty strings which repeats each string infinitely often and let m be large enough and let σ be computed as above. It follows that σ and $R_{T_0}(0)R_{T_0}(1)\dots R_{T_0}(4m)$ differ for an i with $m \leq i < 2m$, that is, satisfy $\sigma(2i+1) \neq R_{T_0}(2i+1)$. Let $w = \sigma(2i+1)$ and $w' = R_{T_0}(2i+1)$. The strings w, w' do not occur in $T_0(0)T_0(1)\dots T_0(m^2)$. Let m' and m'' be least such that $T_0(m') = w$ and $T_0(m'') = w'$. Without loss of generality assume $m' < m''$. Let T_1 and T_2 be obtained from T_0 as follows:

- If $T_0(i) = T_0(m')$ then $T_1(i) = \varepsilon$ else $T_1(i) = T_0(i)$;
- If $T_1(i) = T_1(m'')$ then $T_2(i) = \varepsilon$ else $T_2(i) = T_1(i)$.

Furthermore, the index of T_1 has Kolmogorov complexity bounded by $\log(m'+2)$ and the index of T_2 has Kolmogorov complexity bounded by $2 \log(m''+2)$ up to an additive constant. When considering m (and thus m' and m'') large enough, one can absorb this constant into $\log(m'+2)$ and $\log(m''+2)$ respectively, and thus Kolmogorov complexity of T_1 and T_2 are bounded by $3 \log(m''+2)$. Now R_{T_1} coincides with R_{T_0} below $2m'$, and R_{T_2} coincides with R_{T_1} below $2m''$.

In the various claims above (Claim 4.7, Claim 4.8), when using complexity of T being $6 \log(m+2)$, only the initial portion of text T of length at most m^2 was used. Thus, it was enough to have the complexity of some text T' coinciding with T up to first m^2 elements having a complexity below $6 \log(m+2)$. Hence, T_1 and T_2 satisfy the requirements needed in the claims.

Thus, one can conclude that, when M processes text R_{T_2} , for large enough s , the machine M (after having seen the first $2s+1$ elements of R_{T_2}) does not remove more than $c_5 \cdot \log(s+2)$ symbols from the top of the stack $V_s^{T_2}$. Furthermore, if one now replaces the first $4m+1$ members of R_{T_2} by the corresponding members of σ , then one gets that M on this new text R'_{T_2} has the same convergence behaviour as on R_{T_2} ; however, one text is for L_w while the other one is for $L_{w'}$, thus these are texts for two different languages and so M does not learn at least one of these languages. \square

5. CONCLUSION

The starting point of this research is that automatic functions can be characterised using one-tape Turing machines. More precisely, a function is automatic iff it is computed by a position-faithful one-tape Turing machine in linear time. This is the smallest reasonable linear time complexity class and so the automatic functions turn out to sit at the bottom of the corresponding hierarchy. An open problem is whether the corresponding formalisation using alternating linear time position-faithful one-tape Turing machines also characterises the automatic functions.

Automatic functions have been investigated in learning theory in order to model resource-bounded learners. Due to Pitt's delaying trick [28], unrestricted recursive learners can be bounded heavily in the time that they use without losing learning power. However, automatic learners are not able to learn every learnable class, as their ability to memorise data

is insufficient. Therefore, one might ask whether one can replace an automatic learner by a linear-time learner working on a one-tape Turing machine with a tape of length bounded by the longest datum seen so far plus some additional memory.

These additional memory devices are not restricted in length, though restricted in the amount of access the learner has per cycle: In each cycle the learner runs in time linear in the longest example seen so far, updates the base tape and accesses the additional storage devices only to retrieve or store a linear number of symbols. It is shown that two additional work tapes, two additional stacks or one additional queue give full learning power; furthermore, the learning power of one additional stack is properly intermediate and the learning power of one additional work tape is better than no additional work tape. It is an open problem whether there is a difference in the learning power of one and two additional work tapes.

For some special cases and slightly superlinear computation time, it was possible to show that one additional work tape is enough. The methods of this proof do not generalise to the general case.

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