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# THE MÖBIUS FUNCTION OF THE CONSECUTIVE PATTERN POSET 

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#### Abstract

An occurrence of a consecutive permutation pattern $p$ in a permutation $\pi$ is a segment of consecutive letters of $\pi$ whose values appear in the same order of size as the letters in $p$. The set of all permutations forms a poset with respect to such pattern containment. We compute the Möbius function of intervals in this poset, providing what may be called a complete solution to the problem. For most intervals our results give an immediate answer to the question. In the remaining cases, we give a polynomial time algorithm to compute the Möbius function. In particular, we show that the Möbius function only takes the values $-1,0$ and 1 .


## 1. PRELIMINARIES AND INTRODUCTION

For the poset of classical permutation patterns, the first results about its Möbius function were obtained in [SV]. Further results appear in [ST] and BJJS. The general problem in this case of classical patterns seems quite hard. In contrast, the poset of consecutive pattern containment has a much simpler structure. In this paper we compute the Möbius function of that poset and provide what may be called a complete solution to the problem. In most cases our results give an immediate answer. In the remaining cases, we give a polynomial time recursive algorithm to compute the Möbius function. In particular, we show that the Möbius function only takes the values $-1,0$ and 1 .
Unless otherwise specified, all permutations in this paper are taken to be of the set $[d]=\{1,2, \ldots, d\}$ for some positive integer $d$. We denote by $\mathcal{S}_{d}$ the set of all such permutations for a given $d$. An occurrence of a consecutive pattern $\sigma=a_{1} a_{2} \ldots a_{k}$ in a permutation $\tau=b_{1} b_{2} \ldots b_{n}$ is a subsequence $b_{i} b_{i+1} \ldots b_{i+k-1}$ in $\tau$, whose letters appear in the same order of size as the letters in $\sigma$. As an example, there are three occurrences of the consecutive pattern 231 in the permutation 563724891, namely 563,372 and 891. On the other hand, the permutation 253641 avoids 231 , since it contains no occurrence of that pattern.

Consecutive permutation patterns are special cases of the generalized permutation patterns introduced in [BS], and they are not to be confused with the classical permutation patterns, whose occurrences in a permutation do not have to be contiguous. The enumerative properties of occurrences of various consecutive permutation patterns were first studied systematically in EN, but these results will not concern us, as there seems to be no connection between them and the Möbius function studied here.
The set of all permutations forms a poset $\mathcal{P}$ with respect to consecutive pattern containment. In other words, if $\sigma \in \mathcal{S}_{k}$ and $\tau \in \mathcal{S}_{n}$, then $\sigma \leq \tau$ in $\mathcal{P}$ if $\sigma$ occurs as a consecutive pattern in $\tau$. We write $\sigma<\tau$ if $\sigma \leq \tau$ and $\sigma \neq \tau$. As usual in poset terminology, a permutation $\tau$ covers $\sigma$ (and $\sigma$ is covered by $\tau$ ) if $\sigma<\tau$ and there is

[^0]no permutation $\pi$ such that $\sigma<\pi<\tau$. Note that if $\tau$ covers $\sigma$ then $|\tau|-|\sigma|=1$, where $|\pi|$ is the length of $\pi$. The interval $[x, y]$ in a poset $P$, where $x$ and $y$ are elements of $P$, is defined by $[x, y]=\{z \in P \mid x \leq z \leq y\}$. The rank of an interval $[\sigma, \tau]$ in $\mathcal{P}$ is the difference $|\tau|-|\sigma|$. The rank of an element $\pi$ in $[\sigma, \tau]$ is defined to be the rank of the interval $[\sigma, \pi]$.

A filter in a poset $P$ is a set $S \subseteq P$ such that if $x>y$ and $y \in S$, then $x \in S$. An ideal is a set $S \subseteq P$ such that if $x<y$ and $y \in S$, then $x \in S$. A principal filter is a filter with a single minimal element, and a principal ideal is an ideal with a single maximal element. In each case, the single minimal/maximal element is said to generate the filter/ideal.
In the poset $\mathcal{P}$ consider the interval $[\sigma, \tau]$. Our aim is to compute $\mu(\sigma, \tau)$, where $\mu$ is the Möbius function of the incidence algebra of $\mathcal{P}$. The Möbius function is recursively defined by setting $\mu(x, x)=1$ for all $x$, and, if $x \neq y$,

$$
\begin{equation*}
\mu(x, y)=-\sum_{x \leq z<y} \mu(x, z) \tag{1}
\end{equation*}
$$

In particular, if $x \not \leq y$, then $\mu(x, y)=0$.
It is well known (and follows easily from Philip Hall's Theorem, see St, Proposition 3.8.5]) that the Möbius function of an interval $I$ is equal to the Möbius function of the interval $I^{\prime}$ obtained by "turning $I$ upside-down," that is, the interval $I^{\prime}$ for which $x$ is declared to be less than or equal to $y$ when $y \leq x$ in $I$. Thus, we also have

$$
\begin{equation*}
\mu(x, y)=-\sum_{x<z \leq y} \mu(z, y) . \tag{2}
\end{equation*}
$$

Examples of the two ways to compute the Möbius function, described in Equations (11) and (2), are given in Figure (1) Denoting with $[x, y]$ the interval whose Hasse diagram is depicted in the figure, the label of an element $z$ on the left is $\mu(x, z)$, whereas the label of $z$ on the right is $\mu(z, y)$.

Note also that it follows from either (1) or (2) that the sum of $\mu(x, z)$ over all elements $z$ in an interval $[x, y]$ is 0 , that is

$$
\begin{equation*}
\sum_{x \leq z \leq y} \mu(x, z)=0 \tag{3}
\end{equation*}
$$

This identity will be used frequently.


Figure 1. Computing the Möbius function of an interval from bottom to top (left) and from top to bottom (right).

Definition 1.1. Given a sequence of distinct integers $s=s_{1} s_{2} \ldots s_{d}$, the standard form of $s$ is the permutation $\pi=a_{1} a_{2} \ldots a_{d}$ of $\{1,2, \ldots, d\}$ that is order isomorphic to $s$, that is, whose letters appear in the same order of size as those of $s$.

For example, the standard form of both 4731 and 6842 is 3421 , so all three are order isomorphic. Note that, since an occurrence of a consecutive pattern $\rho$ in a permutation $\pi$ only depends on the relative order of the letters in $\rho$, we may as well talk about the occurrence of an arbitrary sequence of distinct integers as a consecutive pattern in $\pi$. For example, the subsequence 352 in 435261 is an occurrence of the sequence 451 .

Definition 1.2. Suppose $\sigma$ occurs in $\tau=a_{1} a_{2} \ldots a_{n}$. If $a_{i+1}$ is the leftmost letter of $\tau$ involved in any occurrence of $\sigma$ in $\tau$, we say that $\tau$ has a left tail of length $i$ with respect to $\sigma$. Analogously, $\tau$ has a right tail of length $j$ with respect to $\sigma$ if $a_{n-j}$ is the rightmost letter of $\tau$ involved in any occurrence of $\sigma$ in $\tau$. If it is clear from the context what $\sigma$ is, we simply talk about left and right tails of $\tau$.

The following definition is borrowed from the theory of codes.
Definition 1.3. Given a permutation $\tau$, its prefix (resp. suffix) pattern of length $k$ is the permutation of length $k$ order isomorphic to the prefix (resp. suffix) of $\tau$ of length $k$. In case the prefix and suffix patterns of length $k$ of $\tau$ coincide, we say that $\tau$ has a bifix pattern of length $k$.

It is useful to note that in our poset $\mathcal{P}$, each permutation can cover at most two different permutations. Namely, if $\sigma$ is covered by $\tau$ then $\sigma$ clearly must occur as all but the first or all but the last letter of $\tau$. Thus, to obtain a permutation covered by $\tau$ we can only remove the first or last letter of $\tau$, thus yielding at most two different permutations. Moreover, the permutations obtained by removing the first and the last letter, respectively, from $\tau$ are order isomorphic if and only if $\tau$ is monotone, that is, if $\tau$ is either the increasing permutation $123 \ldots n$ or the decreasing permutation $n(n-1) \ldots 21$.

In the case where $\sigma$ occurs precisely once in $\tau$, we show that $\mu(\sigma, \tau)$ depends only on the lengths, $a$ and $b$, of the two tails of $\tau$. More precisely, $\mu(\sigma, \tau)$ is 1 if $a=b \leq 1$, it is -1 if $a=0$ and $b=1$ or vice versa, and 0 otherwise (in which case $\tau$ has a tail of length at least 2).

Our main result, Theorem [3.6, deals with intervals $[\sigma, \tau]$ where $\sigma$ occurs at least twice in $\tau$. This result implies that, as in the case of one occurrence, if $\tau$ has a tail of length at least 2 , then $\mu(\sigma, \tau)=0$. In the remaining cases, where the tails of $\tau$ have length at most 1 , the main result gives a recursive algorithm for computing $\mu(\sigma, \tau)$, by producing, if possible, an element $\mathcal{C}$ in $[\sigma, \tau]$, where $|\mathcal{C}|<|\tau|-2$, such that $\mu(\sigma, \tau)=\mu(\sigma, \mathcal{C})$. This element $\mathcal{C}$, if it exists, must be a bifix pattern of $\tau$, and it must lie below the two elements covered by $\tau$, but not below the element obtained by deleting one letter from each end of $\tau$. If no such element $\mathcal{C}$ exists (which is most often the case), we have $\mu(\sigma, \tau)=0$.

## 2. The case of one occurrence

First we consider the case when $\sigma$ occurs precisely once in $\tau$. In this case the interval $[\sigma, \tau]$ can be described very simply. Let $w_{1} \bar{\sigma} w_{2}$ be the factorization of $\tau$, where the entries of $\bar{\sigma}$ constitute the only occurrence of $\sigma$ in $\tau$ and $\left|w_{1}\right|=a$, $\left|w_{2}\right|=b$, so $\tau$ has left and right tails of lengths $a$ and $b$, respectively. We obtain the permutations covered by $\tau$ by deleting the first or the last entry of $\tau$ and renaming the remaining
entries. Starting from $\tau$ and deleting elements from both tails in such a way, we are eventually left with $\sigma$. It is easy to see that the interval $(\sigma, \tau)$ is a grid, that is, a direct product of two chains of lengths $a$ and $b$, respectively. See Figure 2 for an example. It is also easy to see that the Möbius function of $(\sigma, \tau)$ for such a grid is 0 , except when $a$ and $b$ are both at most 1 . This is recorded in the following theorem.


Figure 2. The interval [123, 68513427]
Theorem 2.1. Suppose $\sigma$ occurs precisely once in $\tau$, and that $\tau$ has tails of lengths $a$ and $b$, respectively. Then $\mu(\sigma, \tau)=1$ if $a=b=0$ or $a=b=1$, and $\mu(\sigma, \tau)=-1$ if $a=0, b=1$, or $a=1, b=0$. Otherwise, $\mu(\sigma, \tau)$ is 0 .

## 3. More than one occurrence

Lemma 3.1. Suppose $\sigma$ occurs in $\tau$ and that $r=|\tau|-|\sigma| \leq 2$. Then, if $r=0$, $\mu(\sigma, \tau)=1$. If $r=1$, we have $\mu(\sigma, \tau)=-1$. If $r=2$, then $\mu(\sigma, \tau)=0$ if $\tau$ (and thus also $\sigma$ ) is a monotone permutation, otherwise $\mu(\sigma, \tau)=1$.

Proof. If $r<2$, the claim is obvious, since $[\sigma, \tau]$ is either a singleton or a chain of two elements. If $r=2$, and $\tau$ is monotone, removing a letter from either end of $\tau$ yields the same (monotone) permutation, so $[\sigma, \tau]$ is a chain of three elements, and $\mu(\sigma, \tau)=0$. If $r=2$ and $\tau$ is not monotone, removing one letter from one end of $\tau$ yields a permutation different from the one obtained by removing a letter from the other end, for otherwise $\tau$ would be monotone. In that case, $[\sigma, \tau]$ is a Boolean algebra with four elements, whose Möbius function is 1.

Because of Lemma3.1, we will from now on only consider intervals $[\sigma, \tau]$ of rank at least three, that is, where $|\tau|-|\sigma| \geq 3$. We also only need to consider pairs $(\sigma, \tau)$ such that $\sigma$ occurs at least twice in $\tau$, since the single-occurrence case is taken care of in Section 2,
We will frequently refer to certain special elements of an interval $[\sigma, \tau]$ described in the following Definition.

Definition 3.2. Given a permutation $\tau$, we let ' $\tau$ be the standard form of $\tau$ after having removed its first letter, $\tau^{\prime}$ be the standard form of $\tau$ after having removed its last letter, and ' $\tau$ ' be the standard form of $\tau$ after having removed both its first and last letter. We refer to ' $\tau$ ' as the interior of $\tau$.

Thus, for instance, if $\tau=68513427$, then $' \tau=7513426, \tau^{\prime}=6751342$ and $^{\prime} \tau^{\prime}=$ 651342.

Lemma 3.3. Suppose $\sigma$ occurs at least twice in $\tau$ and that $|\tau|-|\sigma| \geq 3$. If ' $\tau$ ' does not lie in $[\sigma, \tau]$, we have $\mu(\sigma, \tau)=1$.

Proof. Observe first that $\tau$ cannot be monotone, since then it would be possible to remove one letter from each end and get a permutation in $[\sigma, \tau]$, given that $|\tau|-|\sigma| \geq 3$.
Together with the hypothesis this implies that $\tau$ must have precisely two occurrences of $\sigma$, which necessarily appear at its two ends. Therefore $[\sigma, \tau]$ consists of two chains (of equal lengths) having in common only the minimum $(\sigma)$ and the maximum $(\tau)$. It is easy to see that the Möbius function of such an interval is 1 .

In the lemmas above, and in Section 2 we have taken care of all intervals except those where $\tau$ has at least two occurrences of $\sigma,|\tau|-|\sigma| \geq 3$ and the interior of $\tau$ lies in $[\sigma, \tau]$. We now deal with these remaining cases.
Lemma 3.4. Given $\sigma$ and $\tau$ in $\mathcal{P}$, let

$$
{ }^{\prime} C=\left\{\rho \in[\sigma, \tau] \mid \rho<^{\prime} \tau, \rho \not 又^{\prime} \tau^{\prime}\right\}
$$

and let

$$
C^{\prime}=\left\{\rho \in[\sigma, \tau] \mid \rho<\tau^{\prime}, \rho \not \leq^{\prime} \tau^{\prime}\right\} .
$$

The sets ' $C$ and $C^{\prime}$ are chains, and ${ }^{\prime} C \cap C^{\prime}$ has at most one element. Moreover, if $z \in C^{\prime} \backslash^{\prime} C$ or $z \in^{\prime} C \backslash C^{\prime}$ then $[z, \tau]$ is a chain.

Proof. A permutation $\rho$ in 'C cannot occur in the interior of $\tau$ (since $\rho \not z^{\prime} \tau^{\prime}$ ), and thus it has to be a suffix pattern of $\tau$ (since $\rho<^{\prime} \tau$ ). This implies that ${ }^{\prime} C$ is a chain, since two suffix patterns of $\tau$ must be comparable in $\mathcal{P}$. The argument for $C^{\prime}$ is analogous, so all permutations in $C^{\prime}$ occur as prefix patterns of $\tau$.
Suppose that $x, y \in{ }^{\prime} C \cap C^{\prime}$ and $x \neq y$. We can assume, without loss of generality, that $y<x$, since $x$ and $y$ are elements of the chain $C^{\prime}$. Then $y$ must be a proper prefix pattern of $x$, since both belong to $C^{\prime}$ and thus are prefix patterns of $\tau$. Likewise, since both belong to ${ }^{\prime} C, y$ must be a proper suffix pattern of $x$. But, an element that occurs as a strict suffix of a prefix of $\tau$ and as a strict prefix of a suffix of $\tau$ must occur in the interior of $\tau$, and thus we must have $y \leq \tau^{\prime} \tau^{\prime}$, which contradicts the assumption that $y \in C^{\prime}$. So, ${ }^{\prime} C \cap C^{\prime}$ can contain at most one element.
Assume now that $z \in{ }^{\prime} C \backslash C^{\prime}$. Let $y \in[z, \tau]$ and suppose that $y$ is neither $' \tau$ nor $\tau$ (in which cases $[y, \tau]$ would be a chain). In this case $y$ cannot lie below either $\tau^{\prime}$ or ' $\tau$ ', since then we would have $z<{ }^{\prime} \tau^{\prime}$ (because ' $C$ is a chain), contradicting the hypothesis $z \in{ }^{\prime} C$. Hence, for all $y \in[z, \tau]$, we have $y \in{ }^{\prime} C \cup\left\{{ }^{\prime} \tau, \tau\right\}$, which is a chain. Then $[z, \tau] \subseteq{ }^{\prime} C \cup\left\{^{\prime} \tau, \tau\right\}$, whence $[z, \tau]$ is a chain. An analogous argument shows that $[z, \tau]$ is a chain if $z \in C^{\prime} \backslash^{\prime} C$.

We have thus shown that ${ }^{\prime} C \cap C^{\prime}$ has at most one element. In case it exists, we give it a special name, and record its properties in the following lemma, which is an immediate consequence of the proof of Lemma 3.4

Lemma 3.5 (and Definition). Given $\sigma$ and $\tau$ in $\mathcal{P}$, let ' $C$ and $C^{\prime}$ be as defined in Lemma 3.4. If ' $C \cap C^{\prime}$ is nonempty, let $\mathcal{C}$ be its (necessarily unique) element. Then $\mathcal{C}$ has the following properties:
(1) $\mathcal{C}<{ }^{\prime} \tau$ and $\mathcal{C}<\tau^{\prime}$,
(2) $\mathcal{C} \not \leq^{\prime} \tau^{\prime}$.

Conversely, if $\mathcal{C}$ satisfies the above two conditions, then $\mathcal{C} \in{ }^{\prime} C \cap C^{\prime}$. In what follows, we refer to $\mathcal{C}$ as the carrier element of $[\sigma, \tau]$, or simply as $\mathcal{C}$, if it is clear from the context what $\sigma$ and $\tau$ are.

Observe that $[\sigma, \tau]$ can be expressed as the disjoint union of the principal ideal generated by ' $\tau^{\prime}$ and the filter ' $C \cup C^{\prime} \cup\left\{{ }^{\prime} \tau, \tau^{\prime}, \tau\right\}$. In case $\mathcal{C}$ exists, such a filter is the principal filter generated by $\mathcal{C}$. This fact holds in general, that is, even when the interval $[\sigma, \tau]$ has rank $r \leq 2$. More precisely, if $r=1$, then the ideal $I$ is empty and the filter is $F=\{\sigma, \tau\}$, whereas if $r=2$, then $I=\left\{^{\prime} \tau^{\prime}\right\}=\{\sigma\}$ and $F=\left\{{ }^{\prime} \tau, \tau^{\prime}, \tau\right\}$.

Theorem 3.6. Suppose $\tau$ has at least two occurrences of $\sigma$ and that $|\tau|-|\sigma| \geq$ 3. Assume that ' $\tau$ ' lies in $[\sigma, \tau]$. Then, if $[\sigma, \tau]$ has no carrier element, we have $\mu(\sigma, \tau)=0$. Otherwise, $\mu(\sigma, \tau)=\mu(\sigma, \mathcal{C})$.

Proof. In our argument we will compute the Möbius function of $[\sigma, \tau]$ from top to bottom, using the formula for the Möbius function in Equation (2). We will thus compute $\mu(\sigma, \tau)$ by recursively computing the value of $\mu(z, \tau)$ for $z$ of decreasing ranks, starting with $z=\tau$, as is done on the right hand side in Figure 1
Note that, since $|\tau|-|\sigma| \geq 3$, we have that $\sigma$ lies strictly below all of $\tau, \tau^{\prime},{ }^{\prime} \tau,{ }^{\prime} \tau^{\prime}$, and $\sigma$ is in particular not equal to any one of them.
If $z \in C^{\prime} \backslash^{\prime} C$ then, by Lemma 3.4, the interval $[z, \tau]$ is a chain. Since $z<\tau^{\prime}<\tau$, this chain has rank at least 2 . Thus, $\mu(z, \tau)=0$. An analogous argument shows that $\mu(z, \tau)=0$ when $z \in{ }^{\prime} C \backslash C^{\prime}$. Thus, every element $t$ in $C^{\prime} \cup^{\prime} C$, except $\mathcal{C}$ (if it exists), has $\mu(t, \tau)=0$.
If $\mathcal{C}$ does not exist, then we claim that $\mu(y, \tau)=0$ whenever $y<{ }^{\prime} \tau^{\prime}$, that is whenever $y$ is different from ' $\tau$ ' but belongs to the principal ideal generated by ' $\tau$ '. We claim that a maximal such element $y$ (which must be covered by ' $\tau$ ') lies below precisely four elements $t$ with $\mu(t, \tau) \neq 0$, namely $\tau, \tau^{\prime},{ }^{\prime} \tau, \tau^{\prime}$, and hence has $\mu(y, \tau)=0$. This is because any other element $z$ that $y$ could lie below must belong to ' $C \cup C^{\prime}$, since $z<{ }^{\prime} \tau^{\prime}$ is impossible by virtue of $y$ being maximal. Thus, $\mu(z, \tau)=0$, according to the preceding paragraph. By induction, this now also applies to all other elements lying strictly below ' $\tau$ ', since each of them lies below precisely four elements $t$ with $\mu(t, \tau) \neq 0$, namely, ${ }^{\prime} \tau^{\prime},{ }^{\prime} \tau, \tau^{\prime}$ and $\tau$. In particular, this shows that $\mu(\sigma, \tau)=0$ if $[\sigma, \tau]$ has no carrier element, and, in fact, that then $\mu(z, \tau)=0$ for every $z \in[\sigma, \tau]$ except for $\tau, \tau^{\prime},{ }^{\prime} \tau$ and ${ }^{\prime} \tau^{\prime}$.

Finally, assume that $\mathcal{C}$ exists. We claim that $\mu(\sigma, \tau)=\mu(\sigma, \mathcal{C})$. Indeed, we have

$$
\mu(\sigma, \tau)=-\sum_{\sigma<z \leq \tau} \mu(z, \tau) .
$$

If $z \not \subset \mathcal{C}$, we have seen above that $\mu(z, \tau)=0$, whence

$$
\begin{equation*}
\mu(\sigma, \tau)=-\sum_{\sigma<z \leq \mathcal{C}} \mu(z, \tau) \tag{4}
\end{equation*}
$$

Next we prove that $\mu(z, \tau)=\mu(z, \mathcal{C})$ whenever $z \leq \mathcal{C}$. We proceed by induction on the difference between the rank of $\mathcal{C}$ and the rank of $z$. If $z=\mathcal{C}$, then $\mu(z, \tau)=1$, since the only elements between $\mathcal{C}$ and $\tau$ having nonzero Möbius function are $\tau,{ }^{\prime} \tau$ and $\tau^{\prime}$. This is the base case of the induction. For any $z \leq \mathcal{C}$, using the induction hypothesis on the elements between $z$ and $\mathcal{C}$ and strictly above $z$, we have

$$
\mu(z, \tau)=-\sum_{z<t \leq \tau} \mu(t, \tau)=-\sum_{z<t \leq \mathcal{C}} \mu(t, \tau)=-\sum_{z<t \leq \mathcal{C}} \mu(t, \mathcal{C})=\mu(z, \mathcal{C})
$$

Plugging this into formula (4) we get

$$
\mu(\sigma, \tau)=-\sum_{\sigma<z \leq \mathcal{C}} \mu(z, \mathcal{C})=\mu(\sigma, \mathcal{C})
$$

as desired.

We thus get a recursion, where we find the carrier element of the interval $[\sigma, \mathcal{C}]$, and iterate this until we get an interval that does not have a carrier element. That final interval either has Möbius function zero, or else its rank is at most 2, making it trivial to compute the Möbius function.
Concerning time complexity, a rough analysis shows that, in the worst case, the recursive procedure described by the above theorem is essentially bounded above by $|\tau|^{3}$. Indeed, supposing that $|\sigma|=k$ and $|\tau|=n$, in order to compute the possible carrier element of $[\sigma, \tau]$, one first tries with the prefix of $\tau$ of length $n-2$, and has to check both that it is isomorphic to the suffix of $\tau$ of length $n-2$ and that it is not isomorphic to ' $\tau$ '. The time needed for this task is proportional to $(n-2)^{2}$. In case no carrier element has been found, one has to try with the prefix of length $n-3$; by an analogous argument, the time needed is proportional to $(n-3)^{2}$. In the worst case, (that is, if $[\sigma, \tau]$ has no $\mathcal{C}$ and $\sigma$ occurs at both ends of $\tau$ ), this task has to be performed until we get to $\sigma$ (i.e. the prefix of $\tau$ of length $k$ ). So the total time needed in this worst case is proportional to $\sum_{i=k}^{n-2} i^{2}$, which is proportional to $n^{3}$.
Observe, however, that the worst case is not the most significant one. Indeed, our results imply that it is difficult for the Möbius function of $[\sigma, \tau]$ to be nonzero when $\sigma$ is very long. To be more precise, when $\tau$ has a tail with respect to $\sigma$ of length at least 2, then the Möbius function of $[\sigma, \tau]$ is equal to 0 . Thus if $\sigma$ is not order isomorphic to four specific subwords of $\tau$ (this is a quadratic test), then the Möbius function is 0 . Moreover, observe that the probability that $\sigma$ appears at the beginning of $\tau$, for instance, is equal to $1 / k$ ! (since we have to choose a $k$-subset of an $n$-set, then to arrange its elements in the unique way which produces a word order-isomorphic to $\sigma$, finally we can permute the remaining $n-k$ elements as we like). So, when the length of the pattern $\sigma$ increases, the probability to have $\sigma$ at the beginning of $\tau$ (as well as in any other specific position) rapidly decreases.
In what follows we will frequently refer to the last $\mathcal{C}$ we find in an interval $[\sigma, \tau]$ when finding the carrier element $\mathcal{C}$ of $[\sigma, \tau]$, then the carrier element of $[\sigma, \mathcal{C}]$ and so on until we have come to the last carrier element $\mathcal{C}^{\prime}$ in this sequence (which implies that $\left[\sigma, \mathcal{C}^{\prime}\right]$ does not have a carrier element). In this case, we will refer to $\mathcal{C}^{\prime}$ as the socle of $[\sigma, \tau]$.

Corollary 3.7. Suppose $\sigma$ occurs in $\tau$ but that the first two (or the last two) letters of $\tau$ are not involved in any occurrence of $\sigma$. Then $\mu(\sigma, \tau)=0$.

Proof. If $\sigma$ occurs only once in $\tau$ the claim has already been established, in Theorem 2.1. Assume then that $\sigma$ occurs at least twice in $\tau$. Then $|\tau|-|\sigma| \geq 3$. According to Theorem 3.6, if $[\sigma, \tau]$ has no carrier element $\mathcal{C}$ satisfying the hypotheses of Theorem [3.6, then $\mu(\sigma, \tau)=0$. Otherwise, if there is such a $\mathcal{C}$, then, as a consequence of its definition, $\mathcal{C}$ is order isomorphic to an initial segment of $\tau$, and to a final segment of $\tau$. Thus, there is no occurrence of $\sigma$ in $\mathcal{C}$ involving the first two (or the last two) letters of $\mathcal{C}$. Iterate now the construction of $\mathcal{C}$ for $[\sigma, \mathcal{C}]$ until we obtain the socle $\mathcal{C}^{\prime}$ of $[\sigma, \tau]$. Since $\mathcal{C}^{\prime}$ has a tail of size at least two with respect to $\sigma$ we have $\left|\mathcal{C}^{\prime}\right|-|\sigma| \geq 2$ and there are two possibilities: if $\left|\mathcal{C}^{\prime}\right|-|\sigma|=2$ it follows that $\left[\sigma, \mathcal{C}^{\prime}\right]$
is a three element chain, so $\mu\left(\sigma, \mathcal{C}^{\prime}\right)=0$. Otherwise, $\left|\mathcal{C}^{\prime}\right|-|\sigma| \geq 3$, in which case Theorem 3.6 implies that $\mu\left(\sigma, \mathcal{C}^{\prime}\right)=0$.

Observe that if $\mathcal{C}$ exists then, since it lies in both ${ }^{\prime} C$ and $C^{\prime}, \mathcal{C}$ must be a bifix pattern of $\tau$ of length at least $a+|\sigma|+b$, where $a$ and $b$ are the lengths of the left and right tails (respectively) of $\tau$ with respect to $\sigma$. Also, if $\mathcal{C}$ exists, it does not appear anywhere else in $\tau$ (this just follows from the definition of $\tau$ ). Here are two examples of the $\mathcal{C}$ associated to an interval: For $\sigma=321$ and $\tau=431825976$, we have $\mathcal{C}=321=\sigma$, so $\mu(\sigma, \tau)=\mu(\sigma, \sigma)=1$, whereas if $\sigma=231$ and $\tau=$ 25714893610 we have $\mathcal{C}=245136$ and $\mu(\sigma, \tau)=\mu(\sigma, \mathcal{C})=0$ since the interval [231, 245136] has no carrier element. Observe that in the latter case the initial and final segments of $\tau$ order isomorphic to $\mathcal{C}$ overlap in the letters 48 .

As a consequence of Corollary 3.7, nonzero values of $\mu(\sigma, \tau)$ can occur only when the tails of $\tau$ have length at most 1 . Below we will always assume that $\tau$ satisfies this hypothesis, and we will give a series of partial conditions that facilitate the computation of the Möbius function in several cases.

Denote with $x$ the sum of the lengths of the tails of $\tau$ with respect to $\sigma$. So $x=0$ means that $\tau$ has two occurrences of $\sigma$, one at each end; $x=1$ means that $\tau$ has one occurrence of $\sigma$ at one end and a tail of length 1 at the other end; finally, $x=2$ means that $\tau$ has two tails of length 1 each.

For simplicity, in what follows we will always assume that, when $x=1, \tau$ has an occurrence of $\sigma$ at its right end (and thus a tail of length 1 at its left end). In case $\sigma$ appears at the left end of $\tau$, we simply have to replace each occurrence of the word "suffix" with the word "prefix" in all the following propositions.
Our first result says that it is very difficult for the Möbius function to take the value $(-1)^{x+1}$.

Proposition 3.8. 1. Suppose $x=0$. If $\tau$ does not have a monotone bifix of length $|\sigma|+1$, then $\mu(\sigma, \tau) \neq-1$.
2. Suppose $x=1$. If $\tau$ does not have a bifix of length $|\sigma|+2$ whose suffix of length $|\sigma|+1$ is monotone, then $\mu(\sigma, \tau) \neq 1$.
3. Suppose $x=2$. Then $\mu(\sigma, \tau) \neq-1$.

Proof. Observe that, in general, any permutation of length $\ell$ having a bifix of length $\ell-1$ is necessarily monotone. Suppose $x=0$. If $\mu(\sigma, \tau)=-1$, then (by Theorem 3.6 and Lemma 3.1) if $[\sigma, \tau]$ has a carrier element, then the socle of $[\sigma, \tau]$ must have length $|\sigma|+1$, so it is necessarily monotone (since $\sigma$ is a bifix of it). Suppose $x=1$. If $\mu(\sigma, \tau)=1$, then (by virtue of Theorem 3.6 as well) the socle must have length $|\sigma|+2$, and a direct inspection shows that its suffix of length $|\sigma|+1$ has to be monotone (since $\sigma$ is a bifix of it). Suppose $x=2$. Then necessarily the socle (if it exists) has length at least $|\sigma|+2$, and so $\mu(\sigma, \tau) \neq-1$.

In this direction, a general result that includes all the cases of (but is weaker than) the previous proposition is the following. The proof is easy and is omitted.
Proposition 3.9. If $\tau$ has a non-monotone suffix of length $|\sigma|+x$, then $\mu(\sigma, \tau) \neq$ $(-1)^{x+1}$.

Next we give an easy necessary condition in order to have $\mu(\sigma, \tau)=0$. For this, we first need a definition. A permutation is said to be monotone (reverse) alternating when it is (reverse) alternating and the two permutations induced by its evenindexed elements and odd-indexed elements are both monotone. For instance, the
permutation 342516 is monotone alternating. Alternating permutations have been extensively studied also in recent years, see for instance the survey [S].
Proposition 3.10. 1. Suppose $x=0$ and that $\sigma$ is not the socle of $[\sigma, \tau]$. If $\tau$ has neither a monotone bifix of length $|\sigma|+1$, nor a monotone (reverse) alternating bifix of length $|\sigma|+2$, then $\mu(\sigma, \tau)=0$.
2. Suppose $x=1$. If $\tau$ has neither a bifix of length $|\sigma|+1$ nor a bifix of length $|\sigma|+2$ whose suffix of length $|\sigma|+1$ is monotone, then $\mu(\sigma, \tau)=0$.
3. Suppose $x=2$. If $\tau$ does not have a bifix of length $\sigma \mid+2$, then $\mu(\sigma, \tau)=0$.

Proof. Suppose that $\mu(\sigma, \tau) \neq 0$. Then $[\sigma, \tau]$ necessarily has a carrier element, and we let $\mathcal{C}^{\prime}$ be the socle of $[\sigma, \tau]$. It is clear that, in order to have $\mu(\sigma, \tau) \neq 0, \mathcal{C}^{\prime}$ must be "short", and more precisely $\left|\mathcal{C}^{\prime}\right| \leq|\sigma|+2$.
If $x=0$, we have two distinct cases (since, by hypothesis, $\left|\mathcal{C}^{\prime}\right| \neq|\sigma|$ ). If $\left|\mathcal{C}^{\prime}\right|=|\sigma|+1$, then $\sigma$ is a bifix of $\mathcal{C}^{\prime}$, and so necessarily $\mathcal{C}^{\prime}$ is a monotone bifix of $\tau$. Finally, assume that $\left|\mathcal{C}^{\prime}\right|=|\sigma|+2$. Of course, also in this case $\sigma$ is a bifix of $\mathcal{C}^{\prime}$. Suppose $\sigma=a_{1} \cdots a_{k}$. Our hypothesis implies that, for all $i$, (i) $a_{2 i}<a_{2 i+1}$ if and only if $a_{2 i+2}<a_{2 i+3}$, and (ii) $a_{2 i+1}<a_{2 i+2}$ if and only if $a_{2 i+3}<a_{2 i+4}$. Therefore the descents of $\sigma$ are completely determined by the first three elements $a_{1}, a_{2}$ and $a_{3}$. Moreover, observe that, if we had either $a_{1}<a_{2}<a_{3}$ or $a_{1}>a_{2}>a_{3}$, then $\mathcal{C}^{\prime}$ (and a fortiori $\sigma$ ) would be monotone, and this contradicts the fact that $\left[\sigma, \mathcal{C}^{\prime}\right]$ has no carrier element. Thus we must have either $a_{1}<a_{2}>a_{3}$ or $a_{1}>a_{2}<a_{3}$, which implies that $\mathcal{C}^{\prime}$ is alternating. Finally, there cannot exist any $i$ such that $a_{2 i}<a_{2 i+2}$ and $a_{2 i+2}>a_{2 i+4}$, since the word $a_{1} \cdots a_{k-2}$ is order isomorphic to the word $a_{3} \cdots a_{k}$. The remaining cases can be treated analogously.
If $x=1$, we have only two possibilities. If $\left|\mathcal{C}^{\prime}\right|=|\sigma|+1$, then $\mathcal{C}^{\prime}$ is clearly a bifix of $\tau$ of length $|\sigma|+1$. If $\left|\mathcal{C}^{\prime}\right|=|\sigma|+2$, then Proposition 3.8 tells us that $\mathcal{C}^{\prime}$ is a bifix of $\tau$ whose suffix of length $|\sigma|+1$ is monotone.
If $x=2$, then necessarily $\left|\mathcal{C}^{\prime}\right|=|\sigma|+2$, and $\mathcal{C}^{\prime}$ is of course a bifix of $\tau$.
A general result in this direction is the following.
Proposition 3.11. Denote with $\omega$ the longest bifix of $\tau$ containing $\sigma$ and having length $\leq|\sigma|+2$. Moreover, denote with $\alpha$ (resp. $\beta$ ) the shortest prefix (resp. suffix) pattern of $\tau$ containing the first (resp. last) two occurrences of $\omega$. If $\omega$ exists and $\alpha \neq \beta$, then $\mu(\sigma, \tau)=0$.

Proof. (Sketch) Using an argument similar to that of the previous proposition, suppose that $\mu(\sigma, \tau) \neq 0$, and let $\mathcal{C}^{\prime \prime}$ be the element of $[\sigma, \tau]$ such that the socle of $[\sigma, \tau]$ is the carrier element of the interval $\left[\sigma, \mathcal{C}^{\prime \prime}\right]$. (Thus, $\mathcal{C}^{\prime \prime}$ is the penultimate carrier element found in the iteration culminating in the socle of $[\sigma, \tau])$. Clearly $\mathcal{C}^{\prime \prime}$ has to be a bifix of $\tau$, and it is easy to see that $\mathcal{C}^{\prime \prime}=\alpha=\beta$.

## References

[BS] E. Babson and E. Steingrímsson: Generalized permutation patterns and a classification of the Mahonian statistics, Sém. Lothar. Combin. 44 (2000), Article B44b, 18 pp.
[BJJS] A. Burstein, V. Jelínek, E. Jelínková and E. Steingrímsson: The Möbius function of separable and decomposable permutations, arXiv:1102.1611v1.
[SV] B. E. Sagan and V. Vatter: The Möbius function of a composition poset, J. Algebraic Combin. 24 (2006), 117-136.
[ST] E. Steingrímsson and B. E. Tenner: The Möbius function of the permutation pattern poset, J. Combinatorics 1 (2010), 39-52.
[EN] S. Elizalde and M. Noy: Consecutive patterns in permutations, Adv. Appl. Math. 30 (2003), 110-125.
[St] R. P. Stanley: Enumerative Combinatorics 1, Cambridge University Press, New York, 1997.
[S] R. P. Stanley: A survey of alternating permutations, Contemp. Math., to appear.

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