# MOTIVIC SERRE INVARIANTS MODULO THE SQUARE OF $\mathbb{L}-1$ 

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#### Abstract

Motivic Serre invariants defined by Loeser and Sebag are elements of the Grothendieck ring of varieties modulo $\mathbb{L}-1$. In this paper, we show that we can lift these invariants to modulo the square of $\mathbb{L}-1$ after tensoring the Grothendieck ring with $\mathbb{Q}$ under certain assumptions.


## 1. Introduction

Let $K$ be a complete discrete valuation field with a perfect residue field $k$. For a smooth projective irreducible $K$-variety $X$, Loeser and Sebag 9 defined the motivic Serre invariant $S(X)$. This invariant belongs to the ring $K_{0}\left(\operatorname{Var}_{k}\right) /(\mathbb{L}-1)$, where $K_{0}\left(\operatorname{Var}_{k}\right)$ is the Grothendieck ring of $k$-varieties and $\mathbb{L}:=\left[\mathbb{A}_{k}^{1}\right]$, the class of an affine line in this ring. Let $K_{0}\left(\operatorname{Var}_{k}\right)_{\mathbb{Q}}:=K_{0}\left(\operatorname{Var}_{k}\right) \otimes_{\mathbb{Z}} \mathbb{Q}$. In this paper, we construct, under a certain assumption, an invariant

$$
\tilde{S}(X) \in K_{0}\left(\operatorname{Var}_{k}\right)_{\mathbb{Q}} /(\mathbb{L}-1)^{2}
$$

which coincides with $S(X)$ in $K_{0}\left(\operatorname{Var}_{k}\right)_{\mathbb{Q}} /(\mathbb{L}-1)$.
Remark 1.1. Loeser and Sebag defined the motivic Serre invariant more generally for smooth quasi-compact separated rigid $K$-spaces. For the sake of simplicity, we consider only the case where $X$ is a projective variety.

Let $\mathcal{O}$ be the valuation ring of $K$. The assumption we will make is that the desingularization theorem and the weak factorization theorem hold; their precise statements are as follows:

## Assumption 1.2.

(1) (Desingularization) There exists a regular projective flat $\mathcal{O}$-scheme $\mathcal{X}$ with the generic fiber $\mathcal{X}_{K}:=\mathcal{X} \otimes_{\mathcal{O}} K=X$ such that the special fiber $\mathcal{X}_{k}:=\mathcal{X} \otimes_{\mathcal{O}}$ k is a simple normal crossing divisor in $\mathcal{X}$. (We call such an $\mathcal{X}$ a regular snc model of $X$.)
(2) (Weak factorization) Let $\mathcal{X}$ and $\mathcal{X}^{\prime}$ be regular snc models of $X$. Then there exist finitely many regular snc models of $X$,

$$
\mathcal{X}_{0}=\mathcal{X}, \mathcal{X}_{1}, \ldots, \mathcal{X}_{n}=\mathcal{X}^{\prime}
$$

[^0]such that for every $i$, either the birational map $\mathcal{X}_{i} \rightarrow \mathcal{X}_{i+1}$ is the blowup along a regular center $Z \subset \mathcal{X}_{i+1, k}$ which has normal crossing ${ }^{11}$ with $\mathcal{X}_{i+1, k}$ or its inverse $\mathcal{X}_{i+1} \rightarrow \mathcal{X}_{i}$ has the same description with $\mathcal{X}_{i+1, k}$ replaced with $\mathcal{X}_{i, k}$.

When $X$ has dimension one, this assumption holds, as is well-known. Indeed the above desingularization theorem in this case follows from the desingularization theorem for excellent surfaces by Abhyankar, Hironaka and Lipman (see [8), while the weak factorization follows from the fact that every proper birational morphism of regular integral noetherian schemes of dimension two factors into a sequence of finitely many blowups at closed points. The last fact is well-known in the case of varieties over an algebraically closed field (for instance, [5, V, Cor. 5.4]) and is valid even in our situation as proved in [7, Th. 4.1] in a more general context. Assumption 1.2 holds also when $k$ has characteristic zero. This follows from the recent generalizations to excellent schemes respectively by Temkin [12, 13] and by Abramovich and Temkin [2] of the Hironaka desingularization theorem and the weak factorization theorem of Abramovich, Karu, Matsuki and Włodarczyk [1].

Let $\mathcal{X}$ be a regular snc model of $X$, let $\mathcal{X}_{\mathrm{sm}}$ be its $\mathcal{O}$-smooth locus and let $\mathcal{X}_{\mathrm{sm}, k}:=\mathcal{X}_{\mathrm{sm}} \otimes_{\mathcal{O}} k$. Then $\mathcal{X}_{\mathrm{sm}}$ is a weak Neron model of $X$ in the sense of [3] and by definition,

$$
S(X)=\left[\mathcal{X}_{\mathrm{sm}, k}\right] \in K_{0}\left(\operatorname{Var}_{k}\right) /(\mathbb{L}-1) .
$$

To define our invariant $\tilde{S}(X)$, we also need information on the non-smooth locus of $\mathcal{X}$. Regard $\mathcal{X}_{k}$ as a divisor and write it as $\mathcal{X}_{k}=\sum_{i \in I} a_{i} D_{i}$, where $D_{i}$ are the irreducible components of $\mathcal{X}_{k}$ and $a_{i}$ are the multiplicities of $D_{i}$ in $\mathcal{X}$ respectively. For a subset $H \subset I$, we define

$$
D_{H}^{\circ}:=\bigcap_{h \in H} D_{h} \backslash \bigcup_{i \in I \backslash H} D_{i} .
$$

When $H=\{i\}$, we abbreviate it to $D_{i}^{\circ}$, and when $H=\{i, j\}$, to $D_{i j}^{\circ}$. These locally closed subsets give the stratification

$$
\mathcal{X}_{k}=\bigsqcup_{\emptyset \neq H \subset I} D_{H}^{\circ}
$$

and the stratification

$$
\mathcal{X}_{\mathrm{sm}, k}=\bigcup_{i \in I: a_{i}=1} D_{i}^{\circ} .
$$

From the second stratification, we see that

$$
S(X)=\sum_{i \in I: a_{i}=1}\left[D_{i}^{\circ}\right] \in K_{0}\left(\operatorname{Var}_{k}\right) /(\mathbb{L}-1) .
$$

Loeser and Sebag proved in the paper cited above that this is independent of the model $\mathcal{X}$ and depends only on $X$.

[^1]Definition 1.3. For a regular snc model $\mathcal{X}$ of $X$, we define

$$
\tilde{S}(\mathcal{X}):=\sum_{i \in I: a_{i}=1}\left[D_{i}^{\circ}\right]+\sum_{\substack{\{i, j\} \subset I: \\\left(a_{i}, a_{j}\right)=1}} \frac{1}{a_{i} a_{j}}\left[D_{i j}^{\circ}\right](1-\mathbb{L})
$$

as an element of $K_{0}\left(\operatorname{Var}_{k}\right)_{\mathbb{Q}} /(\mathbb{L}-1)^{2}$. Here $(a, b)$ denotes the greatest common divisor of $a$ and $b$.

Obviously, the two invariants $S(X)$ and $\tilde{S}(\mathcal{X})$ coincide when they are sent to $K_{0}\left(\operatorname{Var}_{k}\right)_{\mathbb{Q}} /(\mathbb{L}-1)$ by the natural maps.

The following is our main theorem:
Theorem 1.4. Let $X$ be a smooth projective $K$-variety. Under Assumption 1.2, the invariant $\tilde{S}(\mathcal{X})$ is independent of the chosen regular snc model $\mathcal{X}$ and depends only on $X$.

The theorem allows us to think of $\tilde{S}(\mathcal{X})$ as an invariant of $X$ and denote it by $\tilde{S}(X)$, which is what was mentioned at the beginning of this Introduction.

## 2. Preparatory reductions

We generalize the invariant $\tilde{S}(\mathcal{X})$ as follows. Let $\mathcal{X}$ be a regular flat $\mathcal{O}$-scheme of finite type such that $\mathcal{X}_{K}$ is smooth and $\mathcal{X}_{k}=\bigcup_{i \in I} D_{i}$ is a simple normal crossing divisor in $\mathcal{X}$. (We no longer suppose that $\mathcal{X}$ or $\mathcal{X}_{K}$ is projective.) For a constructible subset $C \subset \mathcal{X}_{k}$, we define

$$
\tilde{S}(\mathcal{X}, C):=\sum_{\substack{i \in I: \\ a_{i}=1}}\left[D_{i}^{\circ} \cap C\right]+\sum_{\substack{\{i, j\} \subset I: \\\left(a_{i}, a_{j}\right)=1}} \frac{1}{a_{i} a_{j}}\left[D_{i j}^{\circ} \cap C\right](1-\mathbb{L})
$$

as an element of $K_{0}\left(\operatorname{Var}_{k}\right)_{\mathbb{Q}} /(\mathbb{L}-1)^{2}$.
Let $f: \mathcal{Y} \rightarrow \mathcal{X}$ be the blowup along a smooth irreducible center $Z \subset \mathcal{X}_{k}$ which has normal crossings with $\mathcal{X}_{k}$. Then, $\mathcal{Y}$ is an $\mathcal{O}$-scheme satisfying the same conditions as $\mathcal{X}$ does, and we can similarly define $\tilde{S}\left(\mathcal{Y}, C^{\prime}\right)$ for a constructible subset $C^{\prime} \subset \mathcal{Y}_{k}$.

Theorem 1.4 follows from:
Proposition 2.1. Let $\mathcal{X}$ be as above. For any constructible subset $C \subset \mathcal{X}_{k}$, we have

$$
\tilde{S}(\mathcal{X}, C)=\tilde{S}\left(\mathcal{Y}, f^{-1}(C)\right)
$$

Indeed, Theorem 1.4 is a direct consequence of this proposition with $C=\mathcal{X}_{k}$ and Assumption 1.2 .

In what follows, we will prove this proposition. First we will reduce it to the local situation by using:

## Lemma 2.2.

(1) If $C$ is the disjoint union $\bigsqcup_{s=1}^{l} C_{s}$ of constructible subsets $C_{s}$, then

$$
\tilde{S}(\mathcal{X}, C)=\sum_{s=1}^{l} \tilde{S}\left(\mathcal{X}, C_{s}\right)
$$

(2) Let $\mathcal{X}=\bigcup_{\lambda \in \Lambda} U_{\lambda}$ be an open covering. Suppose that for every constructible subset $C \subset \mathcal{X}_{k}$ and for every $\lambda \in \Lambda$,

$$
\tilde{S}\left(\mathcal{X}, C \cap U_{\lambda}\right)=\tilde{S}\left(\mathcal{Y}, f^{-1}\left(C \cap U_{\lambda}\right)\right)
$$

Then, for every constructible subset $C \subset \mathcal{X}_{k}$, we have

$$
\tilde{S}(\mathcal{X}, C)=\tilde{S}\left(\mathcal{Y}, f^{-1}(C)\right)
$$

Proof. The first assertion is obvious. To show the second one, we first claim that there exists a stratification $C=\bigsqcup_{s=0}^{n} C_{s}$ with $C_{s}$ constructible such that each $C_{s}$ is contained in some $U_{\lambda}$. Indeed we can take $C_{0}$ as $C \cap U_{\lambda}$ such that $C$ and $C_{0}$ have equal dimension, then construct $C_{1}$ applying the same procedure to $C \backslash U_{\lambda}$ and so on.

By the assumption, for every $s, \tilde{S}\left(\mathcal{X}, C_{s}\right)=\tilde{S}\left(\mathcal{Y}, f^{-1}\left(C_{s}\right)\right)$. Now, from the first assertion, we get

$$
\tilde{S}(\mathcal{X}, C)=\sum_{s} \tilde{S}\left(\mathcal{X}, C_{s}\right)=\sum_{s} \tilde{S}\left(\mathcal{Y}, f^{-1}\left(C_{s}\right)\right)=\tilde{S}\left(\mathcal{Y}, f^{-1}(C)\right) .
$$

Let $x \in \mathcal{X}_{k}$ be a closed point and take a local coordinate system $x_{1}, \ldots, x_{d} \in$ $\mathcal{O}_{\mathcal{X}, x}$. By shrinking $\mathcal{X}$ if necessary, we may suppose that $x_{1}, \ldots, x_{d}$ are global sections of $\mathcal{O}_{\mathcal{X}}$ and that the special fiber $\mathcal{X}_{k}$ is the zero locus of $\prod_{i=1}^{d^{\prime}} x_{i}, d^{\prime} \leq d$ (thus we identify $I$ with $\left\{1, \ldots, d^{\prime}\right\}$ ) and $Z$ is the common zero locus of $x_{j}, j \in J$, for some subset $J \subset\{1, \ldots, d\}$. From the first assertion of the above lemma, since we obviously have

$$
\tilde{S}(\mathcal{X}, C \backslash Z)=\tilde{S}\left(\mathcal{Y}, f^{-1}(C \backslash Z)\right)
$$

we may also assume that

$$
\begin{equation*}
C \subset Z \tag{2.1}
\end{equation*}
$$

In a few following sections, we will prove Proposition 2.1 in this situation, discussing separately in the cases $(\sharp I=) d^{\prime}=1, d^{\prime}=2$ and $d^{\prime} \geq 3$. Before that, we prepare some notation and a lemma.

Notation 2.3. For $i \in I$, let $D_{i}$ be the prime divisor of $\mathcal{X}$ given by $x_{i}=0$ and let $E_{i} \subset \mathcal{Y}_{k}$ be its strict transform. Let $E_{0} \subset \mathcal{Y}_{k}$ be the exceptional divisor of the blowup $f: \mathcal{Y} \rightarrow \mathcal{X}$. We denote $f^{-1}(C)$ by $\tilde{C}$.

The multiplicity of $E_{i}$ in $\mathcal{Y}_{k}$ is $a_{i}$ for $i \in I$ and

$$
\begin{equation*}
a_{0}:=\sum_{Z \subset D_{i}} a_{i} \tag{2.2}
\end{equation*}
$$

for $i=0$. We will use the following lemma several times.
Lemma 2.4. For $i \in I \backslash J$, if $C \subset Z \cap D_{i}$, then we have $\tilde{C} \subset E_{i}$.
Proof. The morphism $\tilde{C} \rightarrow C$ is a $\mathbb{P}^{\sharp J-1}$-bundle. The divisor $E_{i}$ is the blowup of $D_{i}$ along $Z \cap D_{i}$, which has codimension $\sharp J$ in $D_{i}$. It follows that $E_{i} \cap \tilde{C} \rightarrow C$ is also a $\mathbb{P}^{\sharp J-1}$-bundle. Hence $\tilde{C}$ and $E_{i} \cap \tilde{C}$ coincide and the lemma follows.

$$
\text { 3. The case } d^{\prime}=1
$$

We now begin the proof of Proposition 2.1 in the situation described just before Notation 2.3, In this section, we consider the case $d^{\prime}=1$.

Since $Z \subset \mathcal{X}_{k}$, recalling $I=\left\{1, \ldots, d^{\prime}\right\}$, we see that $1 \in J$. Then

$$
\tilde{S}(\mathcal{X}, C)= \begin{cases}{[C]} & \left(a_{1}=1\right) \\ 0 & (\text { otherwise })\end{cases}
$$

From (2.2), $a_{0}=a_{1}$, and $\left(a_{0}, a_{1}\right)=a_{1}$. Hence, if $a_{1} \neq 1$, then

$$
\tilde{S}(\mathcal{Y}, \tilde{C})=0=\tilde{S}(\mathcal{X}, C)
$$

If $a_{1}=1$, then recalling that $C \subset Z$, we see that $\tilde{C} \subset E_{0}=f^{-1}(Z)$ and that

$$
\tilde{S}(\mathcal{Y}, \tilde{C})=\left[\tilde{C} \backslash E_{1}\right]+\left[E_{1} \cap \tilde{C}\right](1-\mathbb{L})
$$

To compute the right hand side of this equality, we first observe that $\tilde{C}$ is a $\mathbb{P}^{\sharp J-1}{ }_{-}$ bundle over $C$. The divisor $E_{1}$ is the blowup of $D_{1}$ along $Z$. Therefore $E_{1} \cap \tilde{C}$ is a $\mathbb{P}^{\sharp J-2}$-bundle over $C$. Hence

$$
\begin{aligned}
\tilde{S}(\mathcal{Y}, \tilde{C}) & =[C]\left(\left[\mathbb{P}^{\sharp J-1}\right]-\left[\mathbb{P}^{\sharp J-2}\right]\right)+[C]\left[\mathbb{P}^{\sharp J-2}\right](1-\mathbb{L}) \\
& =[C]\left(\mathbb{L}^{\sharp J-1}+\left(1+\mathbb{L}+\cdots+\mathbb{L}^{\sharp J-2}\right)(1-\mathbb{L})\right) \\
& =[C]\left(\mathbb{L}^{\sharp J-1}+1-\mathbb{L}^{\sharp J-1}\right) \\
& =[C] \\
& =\tilde{S}(\mathcal{X}, C) .
\end{aligned}
$$

We conclude that if $d^{\prime}=1$, then $\tilde{S}(\mathcal{X}, C)=\tilde{S}(\mathcal{Y}, \tilde{C})$.

## 4. The case $d^{\prime}=2$

Next we consider the case $d^{\prime}=2$. We have

$$
C=\left(C \cap D_{1}^{\circ}\right) \sqcup\left(C \cap D_{2}^{\circ}\right) \sqcup\left(C \cap D_{12}^{\circ}\right) .
$$

From the case $\sharp I=1$ treated in the last section, we have

$$
\tilde{S}\left(\mathcal{X}, C \cap D_{i}^{\circ}\right)=\tilde{S}\left(\mathcal{Y}, f^{-1}\left(C \cap D_{i}^{\circ}\right)\right) \quad(i=1,2)
$$

Therefore, from Lemma 2.2, replacing $C$ with $C \cap D_{12}^{\circ}$, we may suppose that

$$
\begin{equation*}
C \subset D_{12}^{\circ}=D_{1} \cap D_{2} \tag{4.1}
\end{equation*}
$$

Then we have

$$
\tilde{S}(\mathcal{X}, C)= \begin{cases}\frac{1}{a_{1} a_{2}}[C](1-\mathbb{L}) & \left(\left(a_{1}, a_{2}\right)=1\right) \\ 0 & (\text { otherwise })\end{cases}
$$

We next compute $\tilde{S}(\mathcal{Y}, \tilde{C})$ separately in the case $Z \subset D_{1} \cap D_{2}$ and in the case $Z \not \subset D_{1} \cap D_{2}$.

In the former case, we have $a_{0}=a_{1}+a_{2} \neq 1$ and

$$
\tilde{S}(\mathcal{Y}, \tilde{C})=\sum_{\substack{i \in\{1,2\}: \\\left(a_{0}, a_{i}\right)=1}} \frac{1}{a_{0} a_{i}}\left[\tilde{C} \cap E_{0 i}^{\circ}\right](1-\mathbb{L})
$$

If $\left(a_{1}, a_{2}\right) \neq 1$, then $\left(a_{0}, a_{1}\right) \neq 1$ and $\left(a_{0}, a_{2}\right) \neq 1$, which show that $\tilde{S}(\mathcal{Y}, \tilde{C})=0=$ $\tilde{S}(\mathcal{X}, C)$. If $\left(a_{1}, a_{2}\right)=1$, then we have $\left(a_{0}, a_{1}\right)=\left(a_{0}, a_{2}\right)=1$, and

$$
\tilde{S}(\mathcal{Y}, \tilde{C})=\sum_{i=1}^{2} \frac{1}{a_{0} a_{i}}\left[\tilde{C} \cap E_{0 i}^{\circ}\right](1-\mathbb{L})
$$

Since $E_{1} \cap \tilde{C}=E_{0} \cap E_{1} \cap \tilde{C}_{\tilde{C}} \rightarrow C$ is a trivial $\mathbb{P}^{\sharp J-2}$-bundle and $E_{1} \cap E_{2} \cap \tilde{C} \rightarrow C$ is a hyperplane in it, $E_{01}^{\circ} \cap \tilde{C} \rightarrow C$ is a trivial $\mathbb{A}^{\sharp J-2}$-bundle. (Note that if $\sharp J=2$,
then $E_{1} \cap E_{2}=\emptyset$ and $E_{1} \cap \tilde{C}=E_{01}^{\circ} \cap \tilde{C} \rightarrow C$ is an isomorphism and still a trivial $\mathbb{A}^{\sharp J-2}$-bundle.) Similarly for $E_{02}^{\circ} \cap \tilde{C} \rightarrow C$. Hence

$$
\begin{aligned}
\tilde{S}(\mathcal{Y}, \tilde{C}) & =\left(\frac{1}{\left(a_{1}+a_{2}\right) a_{1}}+\frac{1}{\left(a_{1}+a_{2}\right) a_{2}}\right)[C] \mathbb{L}^{\sharp J-2}(1-\mathbb{L}) \\
& =\frac{1}{a_{1} a_{2}}[C] \mathbb{L}^{\sharp J-2}(1-\mathbb{L}) \\
& \star \frac{1}{=} \frac{1}{a_{1} a_{2}}[C](1-\mathbb{L}) \\
& =\tilde{S}(\mathcal{X}, C) .
\end{aligned}
$$

Here the equality marked with $\star$ follows from

$$
\mathbb{L}(1-\mathbb{L})=(\mathbb{L}-1)(1-\mathbb{L})+1-\mathbb{L}=1-\mathbb{L} \quad \bmod (\mathbb{L}-1)^{2} .
$$

In the case $Z \not \subset D_{1} \cap D_{2}$, we have either $Z \subset D_{1}$ or $Z \subset D_{2}$. Since the two cases are similar, we only discuss the former case. Since $2 \in I \backslash J$, from assumptions (2.1) and (4.1) and Lemma [2.4, we have $\tilde{C} \subset E_{0} \cap E_{2}$. Since $a_{0}=a_{1}, \tilde{C} \rightarrow C$ is a $\mathbb{P}^{\sharp} J-1$-bundle and $\tilde{C} \cap E_{1} \rightarrow C$ is a $\mathbb{P}^{\sharp J-2}$-bundle, we have

$$
\begin{aligned}
\tilde{S}(\mathcal{Y}, \tilde{C}) & =\frac{1}{a_{0} a_{2}}\left[\tilde{C} \cap E_{0,2}^{\circ}\right](1-\mathbb{L}) \\
& =\frac{1}{a_{1} a_{2}}\left[\tilde{C} \backslash E_{1}\right](1-\mathbb{L}) \\
& =\frac{1}{a_{1} a_{2}}[C]\left[\mathbb{P}^{\sharp J-1} \backslash \mathbb{P}^{\sharp J-2}\right](1-\mathbb{L}) \\
& =\frac{1}{a_{1} a_{2}}[C] \mathbb{L}^{\sharp J-1}(1-\mathbb{L}) \\
& =\frac{1}{a_{1} a_{2}}[C](1-\mathbb{L}) \\
& =\tilde{S}(\mathcal{X}, C) .
\end{aligned}
$$

We have completed the proof that $\tilde{S}(\mathcal{Y}, \tilde{C})=\tilde{S}(\mathcal{X}, C)$, when $d^{\prime}=2$.

## 5. The case $d^{\prime} \geq 3$

As in the last section, by induction on $\sharp I$, we may suppose that

$$
\begin{equation*}
C \subset \bigcap_{i \in I} D_{i} . \tag{5.1}
\end{equation*}
$$

Then $\tilde{S}(\mathcal{X}, C)=0$. On the other hand, $\tilde{S}(\mathcal{Y}, \tilde{C})$ is a $\mathbb{Q}$-linear combination of

$$
A_{i}:=\left[\tilde{C} \cap E_{0 i}^{\circ}\right](1-\mathbb{L}), \quad i \in I,
$$

and

$$
B:=\delta_{1, a_{0}}\left[\tilde{C} \cap E_{0}^{\circ}\right],
$$

with $\delta_{1, a_{0}}$ being the Kronecker delta. Thus it suffices to show that $A_{i}=0, i \in I$, and that $B=0$.

We first show that $B=0$. If $\sharp(I \cap J) \geq 2$, then

$$
a_{0}=\sum_{i \in I \cap J} a_{i}>1 .
$$

Hence $B=0$. If $\sharp(I \cap J)<2$, then $I \backslash J$ is non-empty. Assumptions (2.1) and (5.1) and Lemma 2.4 show that $\tilde{C} \cap E_{0}^{\circ}$ is empty, hence $B=0$.

Next we show that $A_{i}=0$. If $\sharp(I \backslash J) \geq 2$, then from Lemma 2.4, for every $i \in I$, there exists $i^{\prime} \in I \backslash\{i\}$ such that $\tilde{C} \subset E_{i^{\prime}}$. Hence $\tilde{C} \cap E_{0 i}^{\circ}=\emptyset$ and $A_{i}=0$.

If $\sharp(I \backslash J)=1$, then by the same reasoning as above, $A_{i}=0$ for $i \in I \cap J$. For $i \in I \backslash J$,

$$
\tilde{C} \cap E_{0 i}^{\circ}=\mathbb{P}_{C}^{\sharp J-1} \backslash \bigcup_{j \in I \cap J} H_{j},
$$

where $\mathbb{P}_{C}^{\sharp J-1}$ denotes the trivial $\mathbb{P}^{\sharp J-1}$-bundle $\mathbb{P}^{\sharp J-1} \times C$ over $C$ and $H_{j}$ are coordinate hyperplanes of $\mathbb{P}_{C}^{\sharp J-1}$. Since $\sharp(I \cap J) \geq 2$,

$$
\begin{aligned}
A_{i} & =[C]\left[\mathbb{G}_{m}^{\sharp(I \cap J)-1} \times \mathbb{A}^{\sharp J-\sharp(I \cap J)}\right](1-\mathbb{L}) \\
& =-[C] \mathbb{L}^{\sharp J-\sharp(I \cap J)}(\mathbb{L}-1)^{\sharp(I \cap J)}=0 \quad \bmod (\mathbb{L}-1)^{2} .
\end{aligned}
$$

If $\sharp(I \backslash J)=0$, equivalently if $Z \subset D_{i}$ for every $i \in I$, then for every $i \in I$,

$$
\tilde{C} \cap E_{0 i}^{\circ}=\mathbb{P}_{C}^{\sharp J-2} \backslash \bigcup_{j \in I \backslash\{i\}} H_{j},
$$

where $H_{j}$ are coordinate hyperplanes of $\mathbb{P}_{C}^{\sharp J-2}$. We have

$$
A_{i}=[C]\left[\mathbb{G}_{m}^{\sharp I-2} \times \mathbb{A}^{\sharp J-\sharp I}\right](1-\mathbb{L})=-[C] \mathbb{L}^{\sharp J-\sharp I}(\mathbb{L}-1)^{\sharp I-1}=0 \quad \bmod (\mathbb{L}-1)^{2} .
$$

We thus have proved that $\tilde{S}(\mathcal{X}, C)=\tilde{S}(\mathcal{Y}, \tilde{C})=0$ also when $d^{\prime} \geq 3$, which completes the proofs of Proposition 2.1 and Theorem 1.4 .

## 6. Closing comments

It is natural to try to refine $\tilde{S}(X)$ further by lifting it to $K_{0}\left(\operatorname{Var}_{k}\right)_{\mathbb{Q}} /(\mathbb{L}-1)^{n}$ for $n>2$ and by adding extra terms of the form

$$
c\left[D_{H}^{\circ}\right](1-\mathbb{L})^{\sharp H-1}
$$

with $c \in \mathbb{Q}, H \subset I, \sharp H \geq 3$. However the author did not manage to find such a refinement.

The original invariant considered by Serre [11] and denoted by $i(X)$ was defined for a $K$-analytic manifold when the residue field $k$ is finite and lives in $\mathbb{Z} /(\sharp k-1)$. There seems to be no counterpart of $\tilde{S}(X)$ in this context, at least in a naive way, because $\mathbb{Z} \otimes_{\mathbb{Z}} \mathbb{Q}=\mathbb{Q}$ is a field and the ideal generated by $(\sharp k-1)^{2}$ in it is the entire field.

The author has no convincing explanation of the meaning of fractional coefficients appearing in the definition of $\tilde{S}(X)$. However, as a possibly related work, we note that Denef and Loeser [4] previously also considered motivic invariants with coefficients in $\mathbb{Q}$.

Nicaise and Sebag [10, Th. 5.4] gave a nice interpretation of the Euler characteristic representation of $S(X)$ in terms of cohomology of the generic fiber (see also [6] for another proof). It would be interesting to look for a similar interpretation of representations of $\tilde{S}(X)$ or $\tilde{S}(X)$ itself.

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[^1]:    ${ }^{1}$ That $Z$ has normal crossings with $\mathcal{X}_{i+1, k}$ means that for every closed point $x \in \mathcal{X}_{i+1, k}$, there exists a regular system of parameters $x_{1}, \ldots, x_{d} \in \mathcal{O}_{\mathcal{X}_{i+1}, x}$ such that in an open neighborhood of $x$, the support of the special fiber $\mathcal{X}_{i+1, k}$ is the zero locus of $\prod_{v \in V} x_{v}$ for some subset $V \subset\{1, \ldots, d\}$ and $Z$ is the common zero locus of $x_{w}, w \in W$, for some $W \subset\{1, \ldots, d\}$.

