

## On common values of $\phi(n)$ and $\sigma(n), I I$

Kevin Ford and Paul Pollack

For each positive-integer valued arithmetic function $f$, let $\mathscr{V}_{f} \subset \mathbb{N}$ denote the image of $f$, and put $\mathscr{V}_{f}(x):=\mathscr{V}_{f} \cap[1, x]$ and $\mathscr{V}_{f}(x):=\# \mathscr{V}_{f}(x)$. Recently Ford, Luca, and Pomerance showed that $\mathscr{V}_{\phi} \cap \mathscr{V}_{\sigma}$ is infinite, where $\phi$ denotes Euler's totient function and $\sigma$ is the usual sum-of-divisors function. Work of Ford shows that $V_{\phi}(x) \asymp V_{\sigma}(x)$ as $x \rightarrow \infty$. Here we prove a result complementary to that of Ford et al. by showing that most $\phi$-values are not $\sigma$-values, and vice versa. More precisely, we prove that, as $x \rightarrow \infty$,

$$
\#\left\{n \leqslant x: n \in \mathscr{V}_{\phi} \cap \mathscr{V}_{\sigma}\right\} \leqslant \frac{V_{\phi}(x)+V_{\sigma}(x)}{(\log \log x)^{1 / 2+o(1)}} .
$$

## 1. Introduction

1A. Summary of results. For each positive-integer valued arithmetic function $f$, let $\mathscr{V}_{f}$ denote the image of $f$, and put $\mathscr{V}_{f}(x):=\mathscr{V}_{f} \cap[1, x]$ and $V_{f}(x):=\# \mathscr{V}_{f}(x)$. In this paper we are primarily concerned with the cases when $f=\phi$, the Euler totient function, and when $f=\sigma$, the usual sum-of-divisors function. When $f=\phi$, the study of the counting function $V_{f}$ goes back to Pillai [1929], and was subsequently taken up in [Erdős 1935; 1945; Erdős and Hall 1973; 1976; Pomerance 1986; Maier and Pomerance 1988; Ford 1998a] (with an announcement in [Ford 1998b]). From the sequence of results obtained in these papers, we mention Erdős's asymptotic formula [1935] for $\log \left(V_{f}(x) / x\right)$, namely

$$
V_{f}(x)=\frac{x}{(\log x)^{1+o(1)}} \quad(x \rightarrow \infty),
$$

[^0]MSC2010: primary 11N37; secondary 11N64, 11A25, 11N36.
Keywords: Euler function, totient, sum of divisors.
and the much more intricate determination of the precise order of magnitude by Ford,

$$
\begin{equation*}
V_{f}(x) \asymp \frac{x}{\log x} \exp \left(C\left(\log _{3} x-\log _{4} x\right)^{2}+D \log _{3} x-\left(D+\frac{1}{2}-2 C\right) \log _{4} x\right) \tag{1-1}
\end{equation*}
$$

Here $\log _{k}$ denotes the $k$-th iterate of the natural logarithm, and the constants $C$ and $D$ are defined as follows: Let

$$
\begin{equation*}
F(z):=\sum_{n=1}^{\infty} a_{n} z^{n}, \quad \text { where } a_{n}=(n+1) \log (n+1)-n \log n-1 \tag{1-2}
\end{equation*}
$$

Since each $a_{n}>0$ and $a_{n} \sim \log n$ as $n \rightarrow \infty$, it follows that $F(z)$ converges to a continuous, strictly increasing function on $(0,1)$, and $F(z) \rightarrow \infty$ as $z \uparrow 1$. Thus, there is a unique real number $\varrho$ for which

$$
\begin{equation*}
F(\varrho)=1 \quad(\varrho=0.542598586098471021959 \ldots) \tag{1-3}
\end{equation*}
$$

In addition, $F^{\prime}$ is strictly increasing, and $F^{\prime}(\varrho)=5.697758 \ldots$ Then

$$
\begin{aligned}
C & =\frac{1}{2|\log \varrho|}=0.817814 \ldots \\
D & =2 C\left(1+\log F^{\prime}(\varrho)-\log (2 C)\right)-\frac{3}{2}=2.176968 \ldots
\end{aligned}
$$

In [Ford 1998a], it is also shown that (1-1) holds for a wide class of $\phi$-like functions, including $f=\sigma$. Consequently, $V_{\phi}(x) \asymp V_{\sigma}(x)$.

Erdős [1959, p. 172] asked if it could be proved that infinitely many natural numbers appear in both $\mathscr{V}_{\phi}$ and $\mathscr{V}_{\sigma}$ (see also [Erdős and Graham 1980]). This question was recently answered by Ford, Luca, and Pomerance [Ford et al. 2010]. Writing $V_{\phi, \sigma}(x)$ for the number of common elements of $\mathscr{V}_{\phi}$ and $\mathscr{V}_{\sigma}$ up to $x$, they proved that

$$
V_{\phi, \sigma}(x) \geqslant \exp \left((\log \log x)^{c}\right)
$$

for some positive constant $c>0$ and all large $x$ (in [Garaev 2011] this is shown for all constants $c>0$ ). This lower bound is probably very far from the truth. For example, if $p$ and $p+2$ form a twin prime pair, then $\phi(p+2)=p+1=\sigma(p)$; a quantitative form of the twin prime conjecture then implies that $V_{\phi, \sigma}(x) \gg x /(\log x)^{2}$. In Part I of this article [Ford and Pollack 2011], we showed that a stronger conjecture of the same type allows for an improvement. Roughly, our result is as follows:
Theorem A. Assume a strong uniform version of Dickson's prime $k$-tuples conjecture. Then as $x \rightarrow \infty$,

$$
V_{\phi, \sigma}(x)=\frac{x}{(\log x)^{1+o(1)}}
$$

Theorem A suggests that $V_{\phi, \sigma}(x)$ is much larger than we might naively expect. This naturally leads one to inquire about what can be proved in the opposite direction; for instance, could it be that a positive proportion of $\phi$-values are also $\sigma$-values?

| $N$ | $V_{\phi}(N)$ | $V_{\sigma}(N)$ | $V_{\phi, \sigma}(N)$ | $\frac{V_{\phi, \sigma}(N)}{V_{\phi}(N)}$ | $\frac{V_{\phi, \sigma}(N)}{V_{\sigma}(N)}$ |
| :--- | ---: | ---: | ---: | :---: | :---: |
| 10000 | 2374 | 2503 | 1368 | 0.5762426 | 0.5465441 |
| 100000 | 20254 | 21399 | 11116 | 0.5488299 | 0.5194635 |
| 1000000 | 180184 | 189511 | 95145 | 0.5280436 | 0.5020553 |
| 10000000 | 1634372 | 1717659 | 841541 | 0.5149017 | 0.4899348 |
| 100000000 | 15037909 | 15784779 | 7570480 | 0.5034264 | 0.4796063 |
| 1000000000 | 139847903 | 146622886 | 69091721 | 0.4940490 | 0.4712206 |

Table 1. Data on $\phi$-values, $\sigma$-values, and common values up to $N=10^{k}$, from $k=5$ to $k=9$.

The numerical data up to $10^{9}$, exhibited in Table 1, suggests that the proportion of common values is decreasing, but the observed rate of decrease is rather slow.

Our principal result is the following estimate, which implies in particular that almost all $\phi$-values are not $\sigma$-values, and vice versa.

Theorem 1.1. As $x \rightarrow \infty$,

$$
V_{\phi, \sigma}(x) \leqslant \frac{V_{\phi}(x)+V_{\sigma}(x)}{(\log \log x)^{1 / 2+o(1)}} .
$$

The proof of Theorem 1.1 relies on the detailed structure theory of totients as developed in [Ford 1998a]. It would be interesting to know the true rate of decay of $V_{\phi, \sigma}(x) / V_{\phi}(x)$.

1B. Sketch. Since the proof of Theorem 1.1 is rather intricate and involves a number of technical estimates, we present a brief outline of the argument in this section.

We start by discarding a sparse set of undesirable $\phi$ and $\sigma$-values. More precisely, we identify (in Lemma 3.2) convenient sets $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$ with the property that almost all $\phi$-values less than or equal to $x$ have all their preimages in $\mathscr{A}_{\phi}$ and almost all $\sigma$-values less than or equal to $x$ have all their preimages in $\mathscr{A}_{\sigma}$. This reduces us to studying how many $\phi$ and $\sigma$-values arise as solutions to the equation

$$
\phi(a)=\sigma\left(a^{\prime}\right), \quad \text { where } \quad a \in \mathscr{A}_{\phi}, a^{\prime} \in \mathscr{A}_{\sigma} .
$$

Note that to show that $V_{\phi, \sigma}(x) / V_{\phi}(x) \rightarrow 0$, we need only count the number of common $\phi-\sigma$-values of this kind, and not the (conceivably much larger) number of pairs $\left(a, a^{\prime}\right) \in \mathscr{A}_{\phi} \times \mathscr{A}_{\sigma}$ corresponding to these values.

What makes the sets $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$ convenient for us? The properties imposed in the definitions of these sets are of two types, anatomical and structural. By anatomical considerations, we mean general considerations of multiplicative structure as commonly appear in elementary number theory (for example, consideration of the
number and size of prime factors). By structural considerations, we mean those depending for their motivation on the fine structure theory of totients developed by Ford [1998a].

Central to our more anatomical considerations is the notion of a normal prime. Hardy and Ramanujan [1917] showed that almost all natural numbers $\leqslant x$ have $\sim \log \log x$ prime factors, and Erdős [1935] showed that the same holds for almost all shifted primes $p-1 \leqslant x$. Moreover, sieve methods imply that if we list the prime factors of $p-1$ on a double-logarithmic scale, then these are typically close to uniformly distributed in $[0, \log \log p]$. Of course, all of this remains true with $p+1$ in place of $p-1$. We assume that the numbers belonging to $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$ have all their prime factors among this set of normal primes.

If we assume that numbers $n$ all of whose prime factors are normal generate "most" $f$-values (for $f \in\{\phi, \sigma\}$ ), we are led to a series of linear inequalities among the (double-logarithmically renormalized) prime factors of $n$. These inequalities are at the heart of the structure theory of totients as developed in [Ford 1998a]. As one illustration of the power of this approach, mapping the $L$ largest prime factors of $n$ (excluding the largest) to a point in $\mathbb{R}^{L}$, the problem of estimating $V_{f}(x)$ reduces to the problem of finding the volume of a certain region of $\mathbb{R}^{L}$, called the fundamental simplex. In broad strokes, this is how one establishes Ford's bound (1-1). We incorporate these linear inequalities into our definitions of $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$. One particular linear combination of renormalized prime factors appearing in the definition of the fundamental simplex is of particular interest to us (see condition (8) in the definition of $\mathscr{A}_{f}$ in Section 3 below); that we can assume this quantity is less than 1 is responsible for the success of our argument.

Suppose now that we have a solution to $\phi(a)=\sigma\left(a^{\prime}\right)$, where $\left(a, a^{\prime}\right) \in \mathscr{A}_{\phi} \times \mathscr{A}_{\sigma}$. We write $a=p_{0} p_{1} p_{2} \cdots$ and $a^{\prime}=q_{0} q_{1} q_{2} \cdots$, where the sequences of $p_{i}$ and $q_{j}$ are nonincreasing. We cut the first of these lists in two places; at the $k$-th prime $p_{k}$ and at the $L$-th prime $p_{L}$. The precise choice of $k$ and $L$ is somewhat technical; one should think of the primes $p_{i}$ larger than $p_{k}$ as the "large" prime divisors of $a$, those smaller than $p_{k}$ but larger than $p_{L}$ as "small", and those smaller than $p_{L}$ as "tiny". The equation $\phi(a)=\sigma\left(a^{\prime}\right)$ can be rewritten in the form

$$
\begin{align*}
& \left(p_{0}-1\right)\left(p_{1}-1\right)\left(p_{2}-1\right) \cdots\left(p_{k-1}-1\right) f d \\
& \quad=\left(q_{0}+1\right)\left(q_{1}+1\right)\left(q_{2}+1\right) \cdots\left(q_{k-1}+1\right) e \tag{1-4}
\end{align*}
$$

where

$$
\begin{equation*}
f:=\phi\left(p_{k} \cdots p_{L-1}\right), \quad d:=\phi\left(p_{L} p_{L+1} \cdots\right), \quad \text { and } \quad e:=\sigma\left(q_{k} q_{k+1} \cdots\right) . \tag{1-5}
\end{equation*}
$$

To see that (1-4) correctly expresses the relation $\phi(a)=\sigma\left(a^{\prime}\right)$, we recall that the primes $p_{1}, \ldots, p_{k}$ are all large, so that by the "anatomical" constraints imposed in the definition of $\mathscr{A}_{\phi}$, each appears to the first power in the prime factorization
of $a$. An analogous statement holds for the primes $q_{1}, \ldots, q_{k}$; this follows from the general principle, established below, that $p_{i} \approx q_{i}$ provided that either side is not too small. There is one respect in which (1-5) may not be quite right: Since $p_{L}$ is tiny, we cannot assume a priori that $p_{L} \neq p_{L-1}$, and so it may be necessary to amend the definition of $d$ somewhat; we ignore this (ultimately minor) difficulty for now.

To complete the argument, we fix $d$ and estimate from above the number of solutions (consisting of $p_{0}, \ldots, p_{k-1}, q_{0}, \ldots, q_{k-1}, e, f$ ) to the relevant equations of the form (1-4); then we sum over $d$. The machinery facilitating these estimates is encoded in Lemma 4.1, which is proved by a delicate, iterative sieve argument of a kind first introduced in [Maier and Pomerance 1988] and developed in [Ford 1998a, §5]. The hypotheses of that lemma include several assumptions about the $p_{i}$ and $q_{j}$, and about $e, f$, and $d$. All of these rather technical hypotheses are, in our situation, consequences of our definitions of $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$; we say more about some of them in a remark following Lemma 4.1.

Notation. Let $P^{+}(n)$ denote the largest prime factor of $n$, understood so that $P^{+}(1)=1$, and let $\Omega(n, U, T)$ denote the total number of prime factors $p$ of $n$ such that $U<p \leqslant T$, counted according to multiplicity. Constants implied by the Landau $O$ and the Vinogradov $\ll$ and $\gg$ symbols are absolute unless otherwise specified. Symbols in boldface type indicate vector quantities.

## 2. Preliminaries

2A. Anatomical tools. We begin with two tools from the standard chest. The first is a form of the upper bound sieve and the second concerns the distribution of smooth numbers.

Lemma 2.1 (see, e.g., [Halberstam and Richert 1974, Theorem 4.2]). Suppose $A_{1}, \ldots, A_{h}$ are positive integers and $B_{1}, \ldots, B_{h}$ are integers such that

$$
E=\prod_{i=1}^{h} A_{i} \prod_{1 \leqslant i<j \leqslant h}\left(A_{i} B_{j}-A_{j} B_{i}\right) \neq 0 .
$$

## Then

$\#\left\{n \leqslant x: A_{i} n+B_{i}\right.$ prime $\left.(1 \leqslant i \leqslant h)\right\} \ll \frac{x}{(\log x)^{h}} \prod_{p \mid E} \frac{1-\frac{v(p)}{p}}{\left(1-\frac{1}{p}\right)^{h}} \ll \frac{x\left(\log _{2}(|E|+2)\right)^{h}}{(\log x)^{h}}$,
where $\nu(p)$ is the number of solutions of the congruence $\prod\left(A_{i} n+B_{i}\right) \equiv 0(\bmod p)$, and the implied constant may depend on $h$.

Let $\Psi(x, y)$ denote the number of $n \leqslant x$ for which $P^{+}(n) \leqslant y$. The following estimate is due to Canfield, Erdős, and Pomerance:

Lemma 2 .2 [Canfield et al. 1983]. Fix $\epsilon>0$. If $2 \leqslant y \leqslant x$ and $u=\frac{\log x}{\log y}$, then

$$
\Psi(x, y)=x / u^{u+o(u)}
$$

for $u \leqslant y^{1-\epsilon}$, as $u \rightarrow \infty$.
The next lemma supplies an estimate for how often $\Omega(n)$ is unusually large; this may be deduced from the theorems in Chapter 0 of [Hall and Tenenbaum 1988].

Lemma 2.3. The number of integers $n \leqslant x$ for which $\Omega(n) \geqslant \alpha \log _{2} x$ is

$$
\ll \alpha \begin{cases}x(\log x)^{-Q(\alpha)} & \text { if } 1<\alpha<2, \\ x(\log x)^{1-\alpha \log 2} \log _{2} x & \text { if } \alpha \geqslant 2,\end{cases}
$$

where $Q(\lambda)=\int_{1}^{\lambda} \log t d t=\lambda \log (\lambda)-\lambda+1$.
In the remainder of this section, we give a precise meaning to the term "normal prime" alluded to in the introduction and draw out some simple consequences. For $S \geqslant 2$, a prime $p$ is said to be $S$-normal if the following two conditions hold for each $f \in\{\phi, \sigma\}$ :

$$
\Omega(f(p), 1, S) \leqslant 2 \log _{2} S,
$$

and, for every pair of real numbers $(U, T)$ with $S \leqslant U<T \leqslant f(p)$, we have

$$
\begin{equation*}
\left|\Omega(f(p), U, T)-\left(\log _{2} T-\log _{2} U\right)\right|<\sqrt{\log _{2} S \log _{2} T} . \tag{2-1}
\end{equation*}
$$

This definition is slightly weaker than the corresponding definition on [Ford 1998a, p. 13], and so the results from that paper remain valid in our context. As a straightforward consequence of the definition, if $p$ is $S$-normal, $f \in\{\phi, \sigma\}$, and $f(p) \geqslant S$, then

$$
\begin{equation*}
\Omega(f(p)) \leqslant 3 \log _{2} f(p) \tag{2-2}
\end{equation*}
$$

The following lemma is a simple consequence of [Ford 1998a, Lemma 2.10] and (1-1):

Lemma 2.4. For each $f \in\{\phi, \sigma\}$, the number of $f$-values less than or equal to $x$ which have a preimage divisible by a prime that is not $S$-normal is

$$
\ll V_{f}(x)\left(\log _{2} x\right)^{5}(\log S)^{-1 / 6} .
$$

We also record the observation that if $p$ is $S$-normal, then $P^{+}(f(p))$ cannot be too much smaller than $p$, on a double-logarithmic scale.

Lemma 2.5. If $5 \leqslant p \leqslant x$ is an $S$-normal prime and $f(p) \geqslant S$, then

$$
\frac{\log _{2} P^{+}(f(p))}{\log _{2} x} \geqslant \frac{\log _{2} p}{\log _{2} x}-\frac{\log _{3} x+\log 4}{\log _{2} x} .
$$

Proof. We have

$$
P^{+}(f(p)) \geqslant f(p)^{\frac{1}{\Omega(f(p))}} \geqslant f(p)^{\frac{1}{3 \log _{2} f(p)}} \geqslant p^{\frac{1}{4 \log _{2} x}} .
$$

The result follows upon taking the double logarithm of both sides.
2B. Structural tools. In this section, we describe more fully some components of the structure theory of totients alluded to in the introduction. Given a natural number $n$, write $n=p_{0}(n) p_{1}(n) p_{2}(n) \cdots$, where $p_{0}(n) \geqslant p_{1}(n) \geqslant p_{2}(n) \geqslant \cdots$ are the primes dividing $n$ (with multiplicity). For a fixed $x$, we put

$$
x_{i}(n ; x)= \begin{cases}\left(\log _{2} p_{i}(n)\right) /\left(\log _{2} x\right) & \text { if } i<\Omega(n) \text { and } p_{i}(n)>2, \\ 0 & \text { if } i \geqslant \Omega(n) \text { or } p_{i}(n)=2 .\end{cases}
$$

Suppose $L \geqslant 2$ is fixed and that $\xi_{i} \geqslant 0$ for $0 \leqslant i \leqslant L-1$. Recall the definition of the $a_{i}$ from (1-2) and let $\mathscr{S}_{L}(\xi)$ be the set of $\left(x_{1}, \ldots, x_{L}\right) \in \mathbb{R}^{L}$ with $0 \leqslant x_{L} \leqslant$ $x_{L-1} \leqslant \cdots \leqslant x_{1} \leqslant 1$ and

$$
\begin{align*}
& \text { ( } \left.I_{0}\right) \quad a_{1} x_{1}+a_{2} x_{2}+\cdots+a_{L} x_{L} \leqslant \xi_{0} \text {, } \\
& a_{1} x_{2}+a_{2} x_{3}+\cdots+a_{L-1} x_{L} \leqslant \xi_{1} x_{1},  \tag{1}\\
& \vdots \quad \vdots \\
& \left(I_{L-2}\right) \quad a_{1} x_{L-1}+a_{2} x_{L} \leqslant \xi_{L-2} x_{L-2} .
\end{align*}
$$

Define $T_{L}(\xi)$ as the volume ( $L$-dimensional Lebesgue measure) of $\mathscr{S}_{L}(\xi)$. For convenience, let $\mathbf{1}=(1,1, \ldots, 1), \mathscr{S}_{L}=\mathscr{S}_{L}(\mathbf{1})$ (the "fundamental simplex"), and let $T_{L}$ be the volume of $\mathscr{S}_{L}$. Let

$$
L_{0}(x):=\left\lfloor 2 C\left(\log _{3} x-\log _{4} x\right)\right\rfloor,
$$

where $C$ is defined as in the introduction. The next lemma allows us to locate the preimages of almost all $f$-values within suitable sets of the form $\mathscr{S}_{L}(\xi)$.

Lemma 2.6 [Ford 1998a, Theorem 15]. Write $L_{0}=L_{0}(x)$. Suppose $0 \leqslant \Psi<L_{0}$, $L=L_{0}-\Psi$, and let

$$
\xi_{i}=\xi_{i}(x)=1+\frac{1}{10\left(L_{0}-i\right)^{3}} \quad(0 \leqslant i \leqslant L-2) .
$$

The number of $f$-values $v \leqslant x$ with a preimage $n$ for which

$$
\left(x_{1}(n ; x), \ldots, x_{L}(n ; x)\right) \notin \mathscr{S}_{L}(\xi) \quad \text { is } \ll V_{f}(x) \exp \left(-\Psi^{2} / 4 C\right) .
$$

For future use, we collect here some further structural lemmas from [Ford 1998a]. The next result, which follows immediately from our (1-1) and Lemma 4.2 of that article, concerns the size of sums of the shape appearing in the definition of inequality $\left(I_{0}\right)$ above.

Lemma 2.7. Suppose that $L \geqslant 2,0<\omega<\frac{1}{10}$, and $x$ is sufficiently large. The number of $f$-values $v \leqslant x$ with a preimage satisfying

$$
a_{1} x_{1}(n ; x)+\cdots+a_{L} x_{L}(n ; x) \geqslant 1+\omega
$$

is $\ll V_{f}(x)\left(\log _{2} x\right)^{5}(\log x)^{-\omega^{2} /\left(150 L^{3} \log L\right)}$.
We will make heavy use of the following (purely geometric) statement about the simplices $\mathscr{S}_{L}(\xi)$, which appears as [Ford 1998a, Lemma 3.10]. Recall from (1-3) that $\varrho=0.542598 \ldots$ denotes the unique real number with $\sum_{n=1}^{\infty} a_{n} \varrho^{n}=1$.
Lemma 2.8. If $\mathbf{x} \in \mathscr{S}_{L}(\boldsymbol{\xi})$ and $\xi_{0}^{L} \xi_{1}^{L-1} \ldots \xi_{L-2}^{2} \leqslant 1.1$, then $x_{j} \leqslant 3 \varrho^{j-i} x_{i}$ when $i<j$, and $x_{j}<3 \varrho^{j}$ for $1 \leqslant j \leqslant L$.

Define $\mathscr{R}_{L}(\xi ; x)$ as the set of integers $n$ with $\Omega(n) \leqslant L$ and

$$
\left(x_{0}(n ; x), x_{1}(n ; x), \ldots, x_{L-1}(n ; x)\right) \in \mathscr{S}_{L}(\xi) .
$$

For $f \in\{\phi, \sigma\}$, put

$$
R_{L}^{(f)}(\xi ; x)=\sum_{n \in \mathscr{R}_{L}(\xi, x)} \frac{1}{f(n)} .
$$

The next lemma, extracted from [Ford 1998a, Lemma 3.12], relates the magnitude of $R_{L}^{(f)}(\xi ; x)$ to the volume of the fundamental simplex $T_{L}$, whenever $\xi$ is suitably close to $\mathbf{1}$. In that article, it plays a crucial role in the proof of the upper-bound aspect of (1-1).
Lemma 2.9. If $1 /\left(1000 k^{3}\right) \leqslant \omega_{L_{0}-k} \leqslant 1 /\left(10 k^{3}\right)$ for $1 \leqslant k \leqslant L_{0}, \xi_{i}=1+\omega_{i}$ for each $i$, and $L \leqslant L_{0}$, then

$$
R_{L}^{(f)}(\xi ; x) \ll\left(\log _{2} x\right)^{L} T_{L}
$$

for both $f=\phi$ and $f=\sigma$.
While only the case $f=\phi$ of Lemma 2.9 appears in the statement of [ibid., Lemma 3.12], the $f=\sigma$ case follows trivially, since $\sigma(n) \geqslant \phi(n)$. In order to apply Lemma 2.9 , we need estimates for the volume $T_{L}$; this is handled by the next lemma, extracted from [ibid., Corollary 3.4].
Lemma 2.10. Assume $1 \leqslant \xi_{i} \leqslant 1.1$ for $0 \leqslant i \leqslant L-2$ and that $\xi_{0}^{L} \xi_{1}^{L-1} \ldots \xi_{L-2}^{2} \asymp 1$. If $L=L_{0}-\Psi>0$, then

$$
\left(\log _{2} x\right)^{L} T_{L}(\xi) \ll Y(x) \exp \left(-\Psi^{2} / 4 C\right)
$$

Here

$$
\begin{equation*}
Y(x):=\exp \left(C\left(\log _{3} x-\log _{4} x\right)^{2}+D \log _{3} x-\left(D+\frac{1}{2}-2 C\right) \log _{4} x\right) . \tag{2-3}
\end{equation*}
$$

We conclude this section with the following technical lemma, which will be needed when we select the sets $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$ in Section 3 .

Lemma 2.11. For $f \in\{\phi, \sigma\}$ and $y \geqslant 20$,

$$
\begin{equation*}
\sum_{\substack{v \in \mathscr{V}_{f} \\ P^{+}(v) \leqslant y}} \frac{1}{v} \ll \frac{\log _{2} y}{\log _{3} y} Y(y), \tag{2-4}
\end{equation*}
$$

where $Y$ is as defined in (2-3). Moreover, for any $b>0$,

$$
\begin{equation*}
Y\left(\exp \left((\log x)^{b}\right)\right)<_{b} Y(x)\left(\frac{\log _{3} x}{\log _{2} x}\right)^{-2 C \log b} . \tag{2-5}
\end{equation*}
$$

Proof. We split the left-hand sum in (2-4) according to whether or not $v \leqslant y^{\log _{2} y}$. The contribution of the large $v$ is $O(1)$ and so is negligible: Indeed, for $t>y^{\log _{2} y}$, we have $\log t / \log y>\log _{2} y$. Thus, by Lemma 2.2, we have $\Psi(t, y) \ll t /(\log t)^{2}$ (say), and the $O(1)$ bound follows by partial summation. We estimate the sum over small $v$ by ignoring the smoothness condition. Put $X=y^{\log _{2} y}$. Since $V_{f}(t) \asymp(t / \log t) Y(t)$, partial summation gives that

$$
\sum_{\substack{v \in \mathscr{Y}_{f} \\ v \leqslant X}} \frac{1}{v} \ll 1+\int_{3}^{X} \frac{Y(t)}{t \log t} d t=(1+o(1)) Y(X) \frac{\log _{2} X}{\log _{3} X}
$$

as $y \rightarrow \infty$. (The last equality follows, for instance, from L'Hôpital's rule.) Since $\log _{2} X / \log _{3} X \sim \log _{2} y / \log _{3} y$ and $Y(X) \sim Y(y)$, we have (2-4). Estimate (2-5) follows from the definition of $Y$ and a direct computation; here it is helpful to note that if we redefine $X:=\exp \left((\log x)^{b}\right)$, then $\log _{3} X=\log _{3} x+\log b$ and $\log _{4} X=\log _{4} x+O_{b}\left(1 / \log _{3} x\right)$.

## 3. Definition of the sets $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$

We continue fleshing out the introductory sketch, giving precise definitions to the preimage sets $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$. Put

$$
\begin{equation*}
L:=\left\lfloor L_{0}(x)-2 \sqrt{\log _{3} x}\right\rfloor, \quad \xi_{i}:=1+\frac{1}{10\left(L_{0}-i\right)^{3}} \quad(1 \leqslant i \leqslant L) . \tag{3-1}
\end{equation*}
$$

The next lemma is a final technical preliminary.
Lemma 3.1. Let $f \in\{\phi, \sigma\}$. The number of $f$-values $v \leqslant x$ with a preimage $n$ for which
(i) $\left(x_{1}(n ; x), \ldots, x_{L}(n ; x)\right) \in \mathscr{S}_{L}(\xi)$ and
(ii) $n$ has fewer than $L+1$ odd prime divisors (counted with multiplicity) is $\ll V_{f}(x) / \log _{2} x$.

Proof. We treat the case when $f=\phi$; the case when $f=\sigma$ requires only small modifications. We can assume that $x / \log x \leqslant n \leqslant 2 x \log _{2} x$, where the last inequality follows from known results on the minimal order of the Euler function. By Lemma 2.3, we can also assume that $\Omega(n) \leqslant 10 \log _{2} x$. Put $p_{i}:=p_{i}(n)$, as defined in Section 2B. Since $\left(x_{1}(n ; x), \ldots, x_{L}(n ; x)\right) \in \mathscr{S}_{L}(\xi)$ by hypothesis, Lemma 2.8 gives that $x_{2}<3 \varrho^{2}<0.9$, and so $p_{2} \leqslant \exp \left((\log x)^{0.9}\right)$. Thus,

$$
n /\left(p_{0} p_{1}\right)=p_{2} p_{3} \cdots \leqslant \exp \left(10\left(\log _{2} x\right)(\log x)^{0.9}\right)=x^{o(1)}
$$

and so $p_{0} \geqslant x^{2 / 5}$ (say) for large $x$. In particular, we can assume that $p_{0}^{2} \nmid n$.
Suppose now that $n$ has exactly $L_{0}-k+1$ odd prime factors, where we fix $k>L_{0}-L$. Then

$$
v=\left(p_{0}-1\right) \phi\left(p_{1} p_{2} \cdots p_{L_{0}-k}\right) 2^{s}
$$

for some integer $s \geqslant 0$. Using the prime number theorem to estimate the number of choices for $p_{0}$ given $p_{1} \cdots p_{L_{0}-k}$ and $2^{s}$, we obtain that the number of $v$ of this form is

$$
\ll \frac{x}{\log x} \sum_{p_{1} \cdots p_{L_{0}-k}} \frac{1}{\phi\left(p_{1} \cdots p_{L_{0}-k}\right)} \sum_{s \geqslant 0} \frac{1}{2^{s}} .
$$

(We use here that $x /\left(\phi\left(p_{1} \cdots p_{L_{0}-k}\right) 2^{s}\right) \gg p_{0} \geqslant x^{2 / 5}$.) The sum over $s$ is $\ll 1$. To handle the remaining sum, we observe that $p_{1} \cdots p_{L_{0}-k}$ belongs to the set $\mathscr{R}_{L_{0}-k}\left(\xi_{k}, x\right)$, where $\xi_{k}:=\left(\xi_{0}, \ldots, \xi_{L_{0}-k-2}\right)$. Thus, the remaining sum is bounded by

$$
R_{L_{0}-k}^{(\phi)}\left(\xi_{k} ; x\right)=\sum_{m \in \mathscr{R}_{L_{0}-k}\left(\xi_{k}, x\right)} \frac{1}{\phi(m)}
$$

So by Lemmas 2.9 and 2.10, both of whose hypotheses are straightforward to verify, $R_{L_{0}-k}^{(\phi)}\left(\xi_{k} ; x\right) \ll\left(\log _{2} x\right)^{L_{0}-k} T_{L_{0}-k} \leqslant\left(\log _{2} x\right)^{L_{0}-k} T_{L_{0}-k}\left(\xi_{k}\right) \ll Y(x) \exp \left(-k^{2} / 4 C\right)$. Collecting our estimates, we obtain a bound of

$$
\ll \frac{x}{\log x} Y(x) \exp \left(-k^{2} / 4 C\right) \ll V_{\phi}(x) \exp \left(-k^{2} / 4 C\right)
$$

Now since $L_{0}-L>2 \sqrt{\log _{3} x}$, summing over $k>L_{0}-L$ gives a final bound which is

$$
\ll V_{\phi}(x) \exp \left(-\left(\log _{3} x\right) / C\right) \ll V_{\phi}(x) / \log _{2} x
$$

as desired.
For the rest of this paper, we fix $\epsilon>0$ and assume that $x \geqslant x_{0}(\epsilon)$. Put

$$
\begin{equation*}
S:=\exp \left(\left(\log _{2} x\right)^{36}\right), \quad \delta:=\sqrt{\frac{\log _{2} S}{\log _{2} x}}, \quad \omega:=\left(\log _{2} x\right)^{-1 / 2+\epsilon / 2} \tag{3-2}
\end{equation*}
$$

For $f \in\{\phi, \sigma\}$, let $\mathscr{A}_{f}$ be the set of $n=p_{0}(n) p_{1}(n) \cdots$ satisfying $f(n) \leqslant x$ and
(0) $n \geqslant x / \log x$,
(1) every squarefull divisor $m$ of $n$ or $f(n)$ satisfies $m \leqslant \log ^{2} x$,
(2) all of the primes $p_{j}(n)$ are $S$-normal,
(3) $\Omega(f(n)) \leqslant 10 \log _{2} x$ and $\Omega(n) \leqslant 10 \log _{2} x$,
(4) if $d \| n$ and $d \geqslant \exp \left(\left(\log _{2} x\right)^{1 / 2}\right)$, then $\Omega(f(d)) \leqslant 10 \log _{2} f(d)$,
(5) $\left(x_{1}(n ; x), \ldots, x_{L}(n ; x)\right) \in \mathscr{I}_{L}(\xi)$,
(6) $n$ has at least $L+1$ odd prime divisors,
(7) $P^{+}\left(f\left(p_{0}\right)\right) \geqslant x^{1 /\left(\log _{2} x\right)}, p_{1}(n)<x^{1 /\left(100 \log _{2} x\right)}$,
(8) $a_{1} x_{1}(n ; x)+\cdots+a_{L} x_{L}(n ; x) \leqslant 1-\omega$.

The following lemma asserts that a generic $f$-value has all of its preimage in $\mathscr{A}_{f}$.
Lemma 3.2. For each $f \in\{\phi, \sigma\}$, the number of $f$-values $\leqslant x$ with a preimage $n \notin \mathscr{A}_{f}$ is

$$
\ll V_{f}(x)\left(\log _{2} x\right)^{-1 / 2+\epsilon}
$$

Remarks. (i) The $\mathscr{A}_{f}$ not only satisfy Lemma 3.2 but do so economically. In fact, from condition (5) and the work of [Ford 1998a, §4], we have that $\# \mathscr{A}_{f} \ll V_{f}(x)$. Thus, on average, an element of $\mathscr{V}_{f}(x)$ has only a bounded number of preimages from $\mathscr{A}_{f}$. So when we turn in Sections 4 and 5 to counting $\phi$-values arising from solutions to $\phi(a)=\sigma\left(a^{\prime}\right)$, with $\left(a, a^{\prime}\right) \in \mathscr{A}_{\phi} \times \mathscr{A}_{\sigma}$, we expect not to be (excessively) overcounting.
(ii) Of the nine conditions defining $\mathscr{A}_{f}$, conditions (0)-(4) are, in the nomenclature of the introduction, purely anatomical, while conditions (5)-(8) depend to some degree on the fine structure theory of [Ford 1998a]. Condition (8) is particularly critical. It is (8) which ensures that the sieve bounds developed in Section 4 result in a nontrivial estimate for $V_{\phi, \sigma}(x)$. Our inability to replace the exponent $\frac{1}{2}$ on $\log _{2} x$ in Lemma 3.2 (or in Theorem 1.1) by a larger number is also rooted in (8).

Proof. It is clear that the number of values of $f(n)$ corresponding to $n$ failing (0) or (1) is $\ll x \log _{2} x / \log x$, which (recalling (1-1)) is permissible for us. By Lemma 2.4 and our choice of $S$, the number of values of $f(n)$ coming from $n$ failing (2) is $\ll V_{f}(x) / \log _{2} x$. The same holds for values coming from $n$ failing (3), by Lemma 2.3.

Suppose now that $n$ fails condition (4). Then $n$ has a unitary divisor $d \geqslant$ $\exp \left(\left(\log _{2} x\right)^{1 / 2}\right)$ with $\Omega(f(d)) \geqslant 10 \log _{2} f(d)$. Put $w:=f(d)$. Then $w \mid f(n)$, and
$f(n) \ll x \log _{2} x$. So if $w \geqslant x^{1 / 2}$, then the number of possibilities for $f(n)$ is

$$
\ll x \log _{2} x \sum_{\substack{w \geqslant x^{1 / 2} \\ \Omega(w) \geqslant 10 \log _{2} w}} \frac{1}{w} \ll \frac{x \log _{2} x}{\log x},
$$

using Lemma 2.3 to estimate sum over $w$. If $w \leqslant x^{1 / 2}$, we observe that $f(n) / w=$ $f(n / d) \in \mathscr{V}_{f}$; hence, with $Y(x)$ defined as in (2-3), the number of corresponding values of $f(n)$ is

$$
\ll \sum_{\substack{\exp \left(\left(\log _{2} x\right)^{1 / 3}\right) \leqslant w \leqslant x^{1 / 2} \\ \Omega(w) \geqslant 10 \log _{2} w}} V_{f}(x / w) \ll \frac{x}{\log x} Y(x) \sum_{\substack{w \geqslant \exp \left(\left(\log _{2} x\right)^{1 / 3}\right) \\ \Omega(w) \geqslant 10 \log _{2} w}} \frac{1}{w} \ll \frac{V_{f}(x)}{\log _{2} x}
$$

By Lemma 2.6, the number of $f$-values with a preimage failing (5) is

$$
\ll \frac{V_{f}(x)}{\log _{2} x}
$$

According to Lemma 3.1, the number of $f$-values with a preimage satisfying (5) but not (6) is also $\ll V_{f}(x) / \log _{2} x$.

Suppose now that $n$ satisfies (0)-(6). In what follows, we write $x_{i}=x_{i}(n ; x)$. From (5), we have $\xi_{0} \geqslant a_{1} x_{1}+a_{2} x_{2} \geqslant\left(a_{1}+a_{2}\right) x_{2}$, and so $x_{2} \leqslant 0.8$. So from (3),

$$
\begin{equation*}
\frac{n}{p_{0}(n) p_{1}(n)}=p_{2}(n) p_{3}(n) \cdots \leqslant \exp \left(10\left(\log _{2} x\right)(\log x)^{0.8}\right)<x^{1 / 100} \tag{3-3}
\end{equation*}
$$

In particular, $p_{0}>x^{1 / 3}+1$ and $f\left(p_{0}\right)>x^{1 / 3}$, so that $v:=f\left(p_{1} p_{2} \cdots\right) \leqslant x^{2 / 3}$. The prime $p_{0}$ satisfies $f\left(p_{0}\right) \leqslant x / v$. For $z$ with $x^{1 / 3}<z \leqslant x$, the number of primes $p_{0}$ with $f\left(p_{0}\right) \leqslant z$ and $P^{+}\left(f\left(p_{0}\right)\right) \leqslant x^{1 / \log _{2} x}$ is (crudely) bounded by $\Psi\left(z, x^{1 / \log _{2} x}\right) \ll z /(\log x)^{2}$, by Lemma 2.2. So the number of values of $f(n)$ coming from $n$ with $P^{+}\left(f\left(p_{0}\right)\right) \leqslant x^{1 / \log _{2} x}$ is

$$
\ll \sum_{\substack{v \leqslant x^{2 / 3} \\ v \in \mathscr{Y}_{f}}} \sum_{\substack{p: f(p) \leqslant x / v \\ P^{+}(f(p)) \leqslant x^{1 / \log _{2} x}}} 1 \ll \frac{x}{(\log x)^{2}} \sum_{\substack{v \leqslant x^{2 / 3} \\ v \in \mathscr{U}_{f}}} \frac{1}{v} \ll \frac{x}{(\log x)^{2-\epsilon}}
$$

To handle the second condition in (7), observe that since $f\left(p_{0}\right) \leqslant x / v$, the prime number theorem (and the bound $v \leqslant x^{2 / 3}$ ) shows that given $v$, the number of possibilities for $p_{0}$ is $\ll x /(v \log x)$. Suppose that $p_{1}(n)>x^{1 /\left(100 \log _{2} x\right)}$. Then $x_{1}=x_{1}(n ; x) \geqslant 0.999$, and we conclude from $\sum_{i \geqslant 1} a_{i} x_{i} \leqslant \xi_{0}$ that either $x_{2} \leqslant \varrho^{3 / 2}$ or $x_{3} \leqslant \varrho^{5 / 2}$. Writing $v_{2}$ for $f\left(p_{2} p_{3} \cdots\right)$ and $v_{3}$ for $f\left(p_{3} p_{4} \cdots\right)$, we see that the number of such $f$-values is

$$
\begin{aligned}
& \ll \frac{x}{\log x} \sum_{p_{1}} \frac{1}{p_{1}}\left(\sum_{P^{+}\left(v_{2}\right)} \frac{1}{v_{2}}+\sum_{p_{2}} \frac{1}{p_{2}} \sum_{\substack{P^{+}\left(v_{3}\right)}} \frac{1}{v_{3}}\right) \\
& \leqslant \exp \left((\log x)^{e^{5 / 2}}\right) \\
& \ll \frac{x}{\log x} \log _{3} x\left(Y(x)\left(\frac{\log _{3} x}{\log _{2} x}\right)^{\left.e^{3 / 2}\right)}+\left(\log _{2} x\right) Y(x)\left(\frac{\log _{3} x}{\log _{2} x}\right)^{3 / 2}\right) \ll \frac{V_{f}(x)}{\left(\log _{2} x\right)^{1 / 2-\epsilon}},
\end{aligned}
$$

using Lemma 2.11 to estimate the sums over $v_{2}$ and $v_{3}$.
Finally, we consider $n$ for which (0)-(7) hold but where condition (8) fails. By Lemma 2.7, we can assume that $a_{1} x_{1}+\cdots+a_{L} x_{L}<1+\omega$, since the number of exceptional $f$-values is

$$
\ll V(x) \exp \left(-\left(\log _{2} x\right)^{\epsilon / 2}\right) \ll \frac{V(x)}{\log _{2} x} .
$$

Thus,

$$
\begin{equation*}
1-\omega<a_{1} x_{1}+\cdots+a_{L} x_{L}<1+\omega \tag{3-4}
\end{equation*}
$$

while by condition $\left(I_{1}\right)$ in the definition of $\mathscr{S}_{L}(\xi), a_{1} x_{2}+\cdots+a_{L-1} x_{L} \leqslant \xi_{1} x_{1}$. We claim that if $J$ is fixed large enough depending on $\epsilon$, then there is some $2 \leqslant j \leqslant J$ with $x_{j} \leqslant \varrho^{j-\epsilon / 3}$. If not, then for large enough $J$,

$$
\xi_{1} x_{1} \geqslant \sum_{j=1}^{J-1} a_{j} x_{j+1} \geqslant \varrho^{1-\epsilon / 3}\left(a_{1} \varrho+a_{2} \varrho^{2}+\cdots+a_{J} \varrho^{J}\right)>\varrho^{1-\epsilon / 4}
$$

Thus, $x_{1} \geqslant \varrho^{1-\epsilon / 4} \xi_{1}^{-1} \geqslant \varrho^{1-\epsilon / 5}$, and so $\xi_{0} \geqslant \varrho^{-\epsilon / 5}\left(a_{1} \varrho+a_{2} \varrho^{2}+\cdots+a_{J} \varrho^{J}\right) \geqslant \varrho^{-\epsilon / 6}$, which is false. This proves the claim. We assume below that $j \in[2, J]$ is chosen as the smallest index with $x_{j} \leqslant \varrho^{j-\epsilon / 3}$; by condition (1), this implies that all of $p_{1}, \ldots, p_{j-1}$ appear to the first power in the prime factorization of $n$.

Now given $x_{2}, \ldots, x_{L}$, we have from (3-4) that $x_{1} \in[\alpha, \alpha+2 \omega]$ for a certain $\alpha$. Thus,

$$
\sum_{p_{1}} \frac{1}{p_{1}} \ll \omega \log _{2} x=\left(\log _{2} x\right)^{1 / 2+\epsilon / 2}
$$

So the number of $f$-values that arise from $n$ satisfying (0)-(7) but failing (8) is

$$
\begin{aligned}
& \ll \frac{x}{\log x} \sum_{j=2}^{J} \sum_{p_{1}, \ldots, p_{j-1}} \frac{1}{p_{1} \cdots p_{j-1}} \sum_{\substack{P^{+}(v) \leqslant \exp \left((\log x)^{e^{j-\epsilon / 3}}\right)}} \frac{1}{v} \\
& \ll \frac{x}{\log x} \sum_{j=2}^{J}\left(\log _{2} x\right)^{j-3 / 2+\epsilon / 2} Y(x)\left(\frac{\log _{3} x}{\log _{2} x}\right)^{-1+j-\epsilon / 3} \ll V_{f}(x)\left(\log _{2} x\right)^{-1 / 2+\epsilon}
\end{aligned}
$$

This completes the proof of Lemma 3.2.

As a corollary of Lemma 3.2, we have that $V_{\phi, \sigma}(x)$ is bounded, up to an additive error of $\ll\left(V_{\phi}(x)+V_{\sigma}(x)\right) /\left(\log _{2} x\right)^{1 / 2-\epsilon}$, by the number of values $\phi(a)$ that appear in solutions to the equation

$$
\phi(a)=\sigma\left(a^{\prime}\right), \quad \text { where } \quad\left(a, a^{\prime}\right) \in \mathscr{A}_{\phi} \times \mathscr{A}_{\sigma}
$$

In Sections 4 and 5, we develop the machinery required to estimate the number of such values. Ultimately, we find that it is smaller than $\left(V_{\phi}(x)+V_{\sigma}(x)\right) /\left(\log _{2} x\right)^{A}$ for any fixed $A$, which immediately gives Theorem 1.1.

## 4. The fundamental sieve estimate

Lemma 4.1. Let $y$ be large, $k \geqslant 1, l \geqslant 0,30 \leqslant S \leqslant v_{k} \leqslant v_{k-1} \leqslant \cdots \leqslant v_{0}=y$, and $u_{j} \leqslant v_{j}$ for $0 \leqslant j \leqslant k-1$. Put $\delta=\sqrt{\log _{2} S / \log _{2} y}, v_{j}=\log _{2} v_{j} / \log _{2} y$, $\mu_{j}=\log _{2} u_{j} / \log _{2} y$. Suppose that $d$ is a natural number for which $P^{+}(d) \leqslant v_{k}$. Moreover, suppose that both of the following hold:
(a) For $2 \leqslant j \leqslant k-1$, either $\left(\mu_{j}, v_{j}\right)=\left(\mu_{j-1}, v_{j-1}\right)$ or $v_{j} \leqslant \mu_{j-1}-2 \delta$. Also, $v_{k} \leqslant \mu_{k-1}-2 \delta$.
(b) For $1 \leqslant j \leqslant k-2$, we have $v_{j}>v_{j+2}$.

The number of solutions of

$$
\begin{equation*}
\left(p_{0}-1\right) \cdots\left(p_{k-1}-1\right) f d=\left(q_{0}+1\right) \cdots\left(q_{k-1}+1\right) e \leqslant y \tag{4-1}
\end{equation*}
$$

in $p_{0}, \ldots, p_{k-1}, q_{0}, \ldots, q_{k-1}, e, f$ satisfying
(i) $p_{i}$ and $q_{i}$ are $S$-normal primes,
(ii) $u_{i} \leqslant P^{+}\left(p_{i}-1\right), P^{+}\left(q_{i}+1\right) \leqslant v_{i}$ for $0 \leqslant i \leqslant k-1$,
(iii) neither $\phi\left(\prod_{i=0}^{k-1} p_{i}\right)$ nor $\sigma\left(\prod_{i=0}^{k-1} q_{i}\right)$ is divisible by $r^{2}$ for a prime $r \geqslant v_{k}$,
(iv) $P^{+}(e f) \leqslant v_{k} ; \Omega(f) \leqslant 4 l \log _{2} v_{k}$,
(v) $p_{0}-1$ has a divisor $\geqslant y^{1 / 2}$ which is composed of primes $>v_{1}$ is

$$
\ll \frac{y}{d}\left(c \log _{2} y\right)^{6 k}(k+1)^{\Omega(d)}\left(\log v_{k}\right)^{8(k+l) \log (k+1)+1}(\log y)^{-2+\sum_{i=1}^{k-1} a_{i} v_{i}+E},
$$

where $E=\delta \sum_{i=2}^{k}(i \log i+i)+2 \sum_{i=1}^{k-1}\left(v_{i}-\mu_{i}\right)$. Here $c$ is an absolute positive constant.

Remarks. Since the lemma statement is very complicated, it may be helpful to elaborate on how it will be applied in Section 5 below. Given $\left(a, a^{\prime}\right) \in \mathscr{A}_{\phi} \times \mathscr{A}_{\sigma}$ satisfying $\phi(a)=\sigma\left(a^{\prime}\right)$, rewrite the corresponding equation in the form (1-4), with $d, e$, and $f$ as in (1-5). (Here $L$ is as in (3-1), and $k$, given more precisely in the next section, satisfies $k \approx L / 2$.) We are concerned with counting the number of
values $\phi(a)$ which arise from such solutions. We partition the solutions according to the value of $d$, which describes the contribution of the "tiny" primes to $\phi(a)$, and by the rough location of the primes $p_{i}$ and $q_{i}$, which we encode in the selection of intervals $\left[u_{i}, v_{i}\right.$ ] (cf. Lemma 2.5). Finally, we apply Lemma 4.1 and sum over both $d$ and the possible selections of intervals; this gives an estimate for the number of $\phi(a)$ which is smaller than $\left(V_{\phi}(x)+V_{\sigma}(x)\right) /\left(\log _{2} x\right)^{A}$, for any fixed $A$.

In our application, conditions (i)-(v) of Lemma 4.1 are either immediate from the definitions, or are readily deduced from the defining properties of $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$. Conditions (a) and (b) are rooted in the observation that while neighboring primes in the prime factorization of $a$ (or $a^{\prime}$ ) may be close together (requiring us to allow $\left.\left[u_{i+1}, v_{i+1}\right]=\left[u_{i}, v_{i}\right]\right)$, the primes $p_{i}(a)$ and $p_{i+2}(a)$ are forced to be far apart on a double-logarithmic scale. Indeed, since $\left(x_{1}(a ; x), \ldots, x_{L}(a ; x)\right) \in \mathscr{S}_{L}(\xi)$, Lemma 2.8 shows that $x_{i+2}<3 \varrho^{2} x_{i}<0.9 x_{i}$.

Proof. We consider separately the prime factors of each shifted prime lying in each interval $\left(v_{i+1}, v_{i}\right]$. For $0 \leqslant j \leqslant k-1$ and $0 \leqslant i \leqslant k$, let

$$
s_{i, j}(n)=\prod_{\substack{p^{a} \|\left(p_{j}-1\right) \\ p \leqslant v_{i}}} p^{a}, \quad s_{i, j}^{\prime}(n)=\prod_{\substack{p^{a} \|\left(q_{j}+1\right) \\ p \leqslant v_{i}}} p^{a}, \quad s_{i}=d f \prod_{j=0}^{k-1} s_{i, j}=e \prod_{j=0}^{k-1} s_{i, j}^{\prime}
$$

Also, for $0 \leqslant j \leqslant k-1$ and $1 \leqslant i \leqslant k$, let

$$
t_{i, j}=\frac{s_{i-1, j}}{s_{i, j}}, \quad t_{i, j}^{\prime}=\frac{s_{i-1, j}^{\prime}}{s_{i, j}^{\prime}}, \quad t_{i}=\prod_{j=0}^{k-1} t_{i, j}=\prod_{j=0}^{k-1} t_{i, j}^{\prime}
$$

For each solution $\mathscr{A}=\left(p_{0}, \ldots, p_{k-1}, f, q_{0}, \ldots, q_{k-1}, e\right)$ of (4-1), let

$$
\begin{aligned}
\sigma_{i}(\mathscr{A}) & =\left\{s_{i} ; s_{i, 0}, \ldots, s_{i, k-1}, f ; s_{i, 0}^{\prime}, \ldots, s_{i, k-1}^{\prime}, e\right\} \\
\tau_{i}(\mathscr{A}) & =\left\{t_{i} ; t_{i, 0}, \ldots, t_{i, k-1}, 1 ; t_{i, 0}^{\prime}, \ldots, t_{i, k-1}^{\prime}, 1\right\}
\end{aligned}
$$

Defining multiplication of $(2 k+l+2)$-tuples component-wise, we have

$$
\begin{equation*}
\sigma_{i-1}(\mathscr{A})=\sigma_{i}(\mathscr{A}) \tau_{i}(\mathscr{A}) \tag{4-2}
\end{equation*}
$$

Let $\mathfrak{S}_{i}$ denote the set of $\sigma_{i}(\mathscr{A})$ arising from solutions $\mathscr{A}$ of (4-1) and $\mathfrak{T}_{i}$ the corresponding set of $\tau_{i}(\mathscr{A})$. By (4-2), the number of solutions of (4-1) satisfying the required conditions is

$$
\begin{equation*}
\left|\mathfrak{S}_{0}\right|=\sum_{\sigma_{1} \in \mathfrak{S}_{1}} \sum_{\substack{\tau_{1} \in \mathfrak{T}_{1} \\ \sigma_{1} \tau_{1} \in \mathfrak{S}_{0}}} 1 \tag{4-3}
\end{equation*}
$$

First, fix $\sigma_{1} \in \mathfrak{S}_{1}$. By assumption (v) in the lemma, $t_{1,0} \geqslant y^{1 / 2}$. Also,

$$
t_{1}=t_{1,0}=t_{1,0}^{\prime} \leqslant y / s_{1}
$$

$t_{1}$ is composed of primes $>v_{1}$, and $s_{1,0} t_{1}+1$ and $s_{1,0}^{\prime} t_{1}-1$ are prime. Write $t_{1}=t_{1}^{\prime} Q$, where $Q=P^{+}\left(t_{1}\right)$. Since $p_{0}$ is an $S$-normal prime, (2-2) gives that

$$
Q \geqslant t_{1}^{1 / \Omega\left(t_{1}\right)} \geqslant t_{1}^{1 / \Omega\left(p_{0}-1\right)} \geqslant y^{1 /\left(2 \Omega\left(p_{0}-1\right)\right)} \geqslant y^{1 /\left(6 \log _{2} y\right)}
$$

Given $t_{1}^{\prime}$, Lemma 2.1 implies that the number of $Q$ is $O\left(y\left(\log _{2} y\right)^{6} /\left(s_{1} t_{1}^{\prime} \log ^{3} y\right)\right)$. Moreover,

$$
\sum \frac{1}{t_{1}^{\prime}} \leqslant \prod_{v_{1}<p \leqslant y}\left(1+\frac{1}{p}+\frac{1}{p^{2}}+\cdots\right) \ll \frac{\log y}{\log v_{1}}=(\log y)^{1-v_{1}}
$$

Consequently, for each $\sigma_{1} \in \mathfrak{S}_{1}$,

$$
\begin{equation*}
\sum_{\substack{\tau_{1} \in \mathfrak{T}_{1} \\ \sigma_{1} \tau_{1} \in \mathfrak{S}_{0}}} 1 \ll \frac{y\left(\log _{2} y\right)^{6}}{s_{1}(\log y)^{2+v_{1}}} \tag{4-4}
\end{equation*}
$$

Next, suppose $2 \leqslant i \leqslant k$. We now apply an iterative procedure: If $v_{i}<v_{i-1}$, we use the identity

$$
\begin{equation*}
\sum_{\sigma_{i-1} \in \mathfrak{S}_{i-1}} \frac{1}{s_{i-1}}=\sum_{\sigma_{i} \in \mathfrak{S}_{i}} \frac{1}{s_{i}} \sum_{\substack{\tau_{i} \in \mathfrak{T}_{i} \\ \sigma_{i} \tau_{i} \in \mathfrak{S}_{i-1}}} \frac{1}{t_{i}} \tag{4-5}
\end{equation*}
$$

If $v_{i}=v_{i-1}$, then (4-5) remains true but contains no information, and in this case we use the alternative identity

$$
\begin{equation*}
\sum_{\sigma_{i-1} \in \mathfrak{S}_{i-1}} \frac{1}{s_{i-1}}=\sum_{\sigma_{i+1} \in \mathfrak{S}_{i+1}} \frac{1}{s_{i+1}} \sum_{\substack{\tau_{i+1} \in \mathfrak{T}_{i+1} \\ \sigma_{i+1} \tau_{i+1} \in \mathfrak{S}_{i-1}}} \frac{1}{t_{i+1}} \tag{4-6}
\end{equation*}
$$

We consider first the simpler case when $v_{i}<v_{i-1}$. Suppose $\sigma_{i} \in \mathfrak{S}_{i}, \tau_{i} \in \mathfrak{T}_{i}$ and $\sigma_{i} \tau_{i} \in \mathfrak{S}_{i-1}$. By assumption (ii), $t_{i}=t_{i, 0} \cdots t_{i, i-1}=t_{i, 0}^{\prime} \cdots t_{i, i-1}^{\prime}$. In addition, $s_{i, i-1} t_{i, i-1}+1=p_{i-1}$ and $s_{i, i-1}^{\prime} t_{i, i-1}^{\prime}-1=q_{i-1}$ are prime. Let $Q:=P^{+}\left(t_{i, i-1}\right)$, $Q^{\prime}:=P^{+}\left(t_{i, i-1}^{\prime}\right), b:=t_{i, i-1} / Q$ and $b^{\prime}:=t_{i, i-1}^{\prime} / Q^{\prime}$.

We consider separately $\mathfrak{T}_{i, 1}$, the set of $\tau_{i}$ with $Q=Q^{\prime}$ and $\mathfrak{T}_{i, 2}$, the set of $\tau_{i}$ with $Q \neq Q^{\prime}$. First,

$$
\Sigma_{1}:=\sum_{\substack{\tau_{i} \in \mathfrak{T}_{i, 1} \\ \sigma_{i} \tau_{i} \in \mathfrak{S}_{i-1}}} \frac{1}{t_{i}} \leqslant \sum_{t} \frac{h(t)}{t} \max _{b, b^{\prime}} \sum_{Q} \frac{1}{Q}
$$

where $h(t)$ denotes the number of solutions of $t_{i, 0} \cdots t_{i, i-2} b=t=t_{i, 0}^{\prime} \cdots t_{i, i-2}^{\prime} b^{\prime}$, and in the sum on $Q, s_{i, i-1} b Q+1$ and $s_{i, i-1}^{\prime} b^{\prime} Q-1$ are prime. By Lemma 2.1, the number of $Q \leqslant z$ is $\ll z(\log z)^{-3}\left(\log _{2} y\right)^{3}$ uniformly in $b, b^{\prime}$. By partial summation,

$$
\sum_{Q \geqslant u_{i-1}} \frac{1}{Q} \ll\left(\log _{2} y\right)^{3}(\log y)^{-2 \mu_{i-1}}
$$

Also, $h(t)$ is at most the number of dual factorizations of $t$ into $i$ factors each, that is, $h(t) \leqslant i^{2 \Omega(t)}$. $\mathrm{By}(2-1), \Omega(t) \leqslant i\left(v_{i-1}-v_{i}+\delta\right) \log _{2} y=: I$. Also, by assumption (iii), $t$ is squarefree. Thus,

$$
\sum_{t} \frac{h(t)}{t} \leqslant \sum_{j \leqslant I} \frac{i^{2 j} H^{j}}{j!}
$$

where

$$
\sum_{v_{i}<p \leqslant v_{i-1}} \frac{1}{p} \leqslant\left(v_{i-1}-v_{i}\right) \log _{2} y+1=: H
$$

By assumption (a), $v_{i-1}-v_{i} \geqslant 2 \delta$, hence $I \leqslant \frac{3}{2} i H \leqslant \frac{3}{4} i^{2} H$. Hence,

$$
\begin{equation*}
\sum_{t} \frac{h(t)}{t} \leqslant\left(\frac{i^{2} H}{I}\right)^{I} \sum_{j \leqslant I} \frac{I^{j}}{j!}<i^{I} \exp (I)=(\log y)^{(i+i \log i)\left(v_{i-1}-v_{i}+\delta\right)} . \tag{4-7}
\end{equation*}
$$

This gives

$$
\Sigma_{1} \ll\left(\log _{2} y\right)^{3}(\log y)^{-2 \mu_{i-1}+(i+i \log i)\left(\nu_{i-1}-v_{i}+\delta\right)} .
$$

For the sum over $\mathfrak{T}_{i, 2}$, set $t_{i}=t Q Q^{\prime}$. Note that $t Q^{\prime}=t_{i, 0} \cdots t_{i, i-2} b$ and $t Q=t_{i, 0}^{\prime} \cdots t_{i, i-2}^{\prime} b^{\prime}$, so $Q \mid t_{i, 0}^{\prime} \cdots t_{i, i-2}^{\prime} b^{\prime}$ and $Q^{\prime} \mid t_{i, 0} \cdots t_{i, i-2} b$. If we fix the factors divisible by $Q$ and by $Q^{\prime}$, then the number of possible ways to form $t$ is $\leqslant i^{2 \Omega(t)}$ as before. Then

$$
\Sigma_{2}:=\sum_{\substack{\tau_{i} \in \mathfrak{T}_{i, 2}, 2 \\ \sigma_{i} t_{i} \in \mathfrak{S}_{i-1}}} \frac{1}{t_{i}} \leqslant \sum_{t} \frac{i^{2 \Omega(t)+2}}{t} \max _{b, b^{\prime}} \sum_{Q, Q^{\prime}} \frac{1}{Q Q^{\prime}},
$$

where $s_{i, i-1} b Q+1$ and $s_{i, i-1}^{\prime} b^{\prime} Q^{\prime}-1$ are prime. By Lemma 2.1, the number of $Q \leqslant z$ (respectively $\left.Q^{\prime} \leqslant z\right)$ is $\ll z(\log z)^{-2}\left(\log _{2} y\right)^{2}$. Hence,

$$
\sum_{Q, Q^{\prime}} \frac{1}{Q Q^{\prime}} \ll\left(\log _{2} y\right)^{4}(\log y)^{-2 \mu_{i-1}}
$$

Combined with (4-7), this gives $\Sigma_{2} \ll i^{2}\left(\log _{2} y\right)^{4}(\log y)^{-2 \mu_{i-1}+(i+i \log i)\left(v_{i-1}-v_{i}+\delta\right)}$. From (a) and (b), $i^{2} \leqslant k^{2} \leqslant\left(\log _{2} y\right)^{2}$. Adding $\Sigma_{1}$ and $\Sigma_{2}$ shows that for each $\sigma_{i}$,

$$
\begin{equation*}
\sum_{\substack{\tau_{i} \in \mathcal{T}_{i} \\ \sigma_{i} \tau_{i} \in \mathfrak{S}_{i-1}}} \frac{1}{t_{i}} \ll\left(\log _{2} y\right)^{6}(\log y)^{-2 \mu_{i-1}+(i \log i+i)\left(v_{i-1}-v_{i}+\delta\right)} . \tag{4-8}
\end{equation*}
$$

We consider now the case when $v_{i}=v_{i-1}$. Set $Q_{1}:=P^{+}\left(t_{i+1, i-1}\right), Q_{2}:=$ $P^{+}\left(t_{i+1, i}\right), Q_{3}:=P^{+}\left(t_{i+1, i-1}^{\prime}\right)$, and $Q_{4}:=P^{+}\left(t_{i+1, i}^{\prime}\right)$. From (iii), we have that $Q_{1} \neq Q_{2}$ and $Q_{3} \neq Q_{4}$. Moreover, letting $b_{i}$ denote the cofactor of $Q_{i}$ in each
case, we have that

$$
\begin{align*}
s_{i+1, i-1} b_{1} Q_{1}+1 & =p_{i-1}, & s_{i+1, i-1}^{\prime} b_{3} Q_{3}-1 & =q_{i-1}  \tag{4-9}\\
s_{i+1, i} b_{2} Q_{2}+1 & =p_{i}, & s_{i+1, i}^{\prime} b_{4} Q_{4}-1 & =q_{i}
\end{align*}
$$

Since there are now several ways in which the various $Q_{i}$ may coincide, the combinatorics is more complicated than in the case when $v_{i}<v_{i-1}$. We index the cases by fixing the incidence matrix $\left(\delta_{i j}\right)$ with $\delta_{i j}=1$ if $Q_{i}=Q_{j}$ and $\delta_{i j}=0$ otherwise.

Write $D=\operatorname{gcd}\left(Q_{1} Q_{2}, Q_{3} Q_{4}\right)$, and let $Q:=Q_{1} Q_{2} / D$ and $Q^{\prime}:=Q_{3} Q_{4} / D$, so that $D, Q$, and $Q^{\prime}$ are formally determined by $\left(\delta_{i j}\right)$. Then $Q Q^{\prime} \mid t_{i+1}$, and writing $t_{i+1} / D=t Q Q^{\prime}$, we have

$$
\begin{align*}
t Q & =t_{i+1,0} t_{i+1,1} \cdots t_{i+1, i-2} b_{3} b_{4}  \tag{4-10}\\
t Q^{\prime} & =t_{i+1,0}^{\prime} t_{i+1,1}^{\prime} \cdots t_{i+1, i-2}^{\prime} b_{1} b_{2} \tag{4-11}
\end{align*}
$$

We now choose which terms on the right-hand sides of (4-10) and (4-11) contain the prime factors of $Q$ and $Q^{\prime}$, respectively; since $\Omega(Q) \leqslant 2$ and $\Omega\left(Q^{\prime}\right) \leqslant 2$, this can be done in at most $(i+1)^{4}$ ways. Having made this choice, the number of ways to form $t$ is bounded by $(i+1)^{2 \Omega(t)}$, and so

$$
\begin{equation*}
\sum_{\substack{\tau_{i+1} \in \mathfrak{T}_{i+1} \\ \sigma_{i+1} \tau_{i+1} \in \mathfrak{S}_{i-1}}} \frac{1}{t_{i+1}} \leqslant \sum_{t} \frac{(i+1)^{2 \Omega(t)+4}}{t} \max _{b_{1}, b_{2}, b_{3}, b_{4}} \sum \frac{1}{D Q Q^{\prime}} \tag{4-12}
\end{equation*}
$$

It is easy to check that $D Q Q^{\prime}=\prod_{j \in \mathscr{J}} Q_{j}$, where $\mathscr{J}$ indexes the distinct $Q_{j}$. For each $j \in \mathscr{J}$, let $n_{j}$ be the number of linear forms appearing in (4-9) involving $Q_{j}$. Since each of these $n_{j}$ linear forms in $Q_{j}$ is prime, as is $Q_{j}$ itself, Lemma 2.1 implies that the number of possibilities for $Q_{j} \leqslant z$ is $\ll z(\log z)^{-n_{j}-1}\left(\log _{2} y\right)^{n_{j}+1}$, and so

$$
\sum_{Q_{j} \geqslant u_{i-1}} \frac{1}{Q_{j}} \ll\left(\log _{2} y\right)^{n_{j}+1}\left(\log u_{i-1}\right)^{-n_{j}} \ll\left(\log _{2} y\right)^{n_{j}+1}(\log y)^{-n_{j} \mu_{i-1}}
$$

uniformly in the choice of the $b$ 's. Since $\sum_{j \in \mathscr{J}} n_{j}=4$ and $\sum_{j \in \mathscr{J}} 1 \leqslant 4$,

$$
\begin{equation*}
\sum \frac{1}{D Q Q^{\prime}} \leqslant \prod_{j \in \mathscr{J}}\left(\sum_{Q_{j} \geqslant u_{i-1}} \frac{1}{Q_{j}}\right) \ll\left(\log _{2} y\right)^{8}(\log y)^{-4 \mu_{i-1}} \tag{4-13}
\end{equation*}
$$

The calculation (4-7), with $i$ replaced by $i+1$, shows that

$$
\begin{equation*}
\sum_{t} \frac{(i+1)^{2 \Omega(t)}}{t} \leqslant(\log y)^{((i+1)+(i+1) \log (i+1))\left(v_{i}-v_{i+1}+\delta\right)} \tag{4-14}
\end{equation*}
$$

Combining (4-12), (4-13), and (4-14) shows that

$$
\sum_{\substack{\tau_{i+1} \in \mathfrak{T}_{i+1} \\ \sigma_{i+1} \tau_{i+1} \in \mathfrak{S}_{i-1}}} \frac{1}{t_{i+1}} \leqslant(i+1)^{4}\left(\log _{2} y\right)^{8}(\log y)^{-4 \mu_{i-1}+((i+1)+(i+1) \log (i+1))\left(v_{i}-v_{i+1}+\delta\right)}
$$

$$
\leqslant\left(\log _{2} y\right)^{12}(\log y)^{-2 \mu_{i-1}+(i \log i+i)\left(v_{i}-\nu_{i-1}\right)-2 \mu_{i}+((i+1) \log (i+1)+(i+1))\left(v_{i+1}-v_{i}+\delta\right)},
$$

where in the last line we use that $v_{i-1}=v_{i}$ and $(i+1)^{4} \leqslant k^{4} \leqslant\left(\log _{2} y\right)^{4}$.
Using (4-3), (4-5), and (4-6) together with the inequalities (4-4), (4-8), and (4-15), we find that the number of solutions of (4-1) is

$$
\ll y\left(c \log _{2} y\right)^{6 k}(\log y)^{-2-v_{1}+\sum_{i=2}^{k}\left(v_{i-1}-v_{i}+\delta\right)(i \log i+i)-2 \mu_{i-1}} \sum_{\sigma_{k} \in \mathfrak{S}_{k}} \frac{1}{s_{k}},
$$

where $c$ is some positive constant. Note that the exponent of $\log y$ is

$$
\leqslant-2+\sum_{i=1}^{k-1} a_{i} v_{i}+E .
$$

It remains to treat the sum on $\sigma_{k}$. Given $s_{k}^{\prime}=s_{k} / d$, the number of possible $\sigma_{k}$ is at most the number of factorizations of $s_{k}^{\prime}$ into $k+1$ factors times the number of factorizations of $d s_{k}^{\prime}$ into $k+1$ factors, which is at most $(k+1)^{\Omega\left(d s_{k}^{\prime}\right)}(k+1)^{\Omega\left(s_{k}^{\prime}\right)}$. By assumptions (i) and (iv), $\Omega\left(s_{k}^{\prime}\right) \leqslant 4(k+l) \log _{2} v_{k}$. Thus,

$$
\begin{aligned}
\sum_{\sigma_{k} \in \mathfrak{S}_{k}} \frac{1}{s_{k}} & \leqslant \frac{(k+1)^{\Omega(d)}(k+1)^{8(k+l) \log _{2} v_{k}}}{d} \sum_{P+\left(s_{k}^{\prime}\right) \leqslant v_{k}} \frac{1}{s_{k}^{\prime}} \\
& \ll \frac{(k+1)^{\Omega(d)}\left(\log v_{k}\right)^{8(k+l) \log (k+1)+1}}{d}
\end{aligned}
$$

## 5. Counting common values: Application of Lemma 4.1

In this section we prove the following proposition, which combined with Lemma 3.2 immediately yields Theorem 1.1. Throughout the rest of this paper, we adopt the definitions of $L$, the $\xi_{i}, S, \delta$, and $\omega$ from (3-1) and (3-2).

Proposition 5.1. Fix $A>0$. For large $x$, the number of distinct values of $\phi$ (a) that arise from solutions to the equation

$$
\phi(a)=\sigma\left(a^{\prime}\right), \quad \text { with } \quad\left(a, a^{\prime}\right) \in \mathscr{A}_{\phi} \times \mathscr{A}_{\sigma},
$$

is smaller than $\left(V_{\phi}(x)+V_{\sigma}(x)\right) /\left(\log _{2} x\right)^{A}$.
Let us once again recall the strategy outlined in the introduction and in the remarks following Lemma 4.1. Let $\left(a, a^{\prime}\right) \in \mathscr{A}_{\phi} \times \mathscr{A}_{\sigma}$ be a solution to $\phi(a)=\sigma\left(a^{\prime}\right)$. Let $p_{i}:=p_{i}(a)$ and $q_{i}:=p_{i}\left(a^{\prime}\right)$, in the notation of Section 2B. We choose a cutoff
$k$ so that all of $p_{0}, \ldots, p_{k-1}$ and $q_{0}, \ldots, q_{k-1}$ are "large". Then by condition (1) in the definition of the sets $\mathscr{A}_{f}$, neither $p_{i}^{2} \mid a$ nor $q_{i}^{2} \mid a^{\prime}$, for $0 \leqslant i \leqslant k-1$. Fixing a notion of "small" and "tiny", we rewrite the equation $\phi(a)=\sigma\left(a^{\prime}\right)$ in the form

$$
\begin{equation*}
\left(p_{0}-1\right) \cdots\left(p_{k-1}-1\right) f d=\left(q_{0}+1\right) \cdots\left(q_{k-1}+1\right) e, \tag{5-1}
\end{equation*}
$$

where $f$ is the contribution to $\phi(a)$ from the "small" primes, $d$ is the contribution from the "tiny" primes, and $e$ is the contribution of both the "small" and "tiny" primes to $\sigma\left(a^{\prime}\right)$.

We then fix $d$ and numbers $u_{i}$ and $v_{i}$, chosen so that $u_{i} \leqslant P^{+}\left(p_{i}-1\right), P^{+}\left(q_{i}+\right.$ $1) \leqslant v_{i}$ for each $0 \leqslant i \leqslant k-1$. With these fixed, Lemma 4.1 provides us with an upper bound on the number of corresponding solutions to (5-1). Such a solution determines the common value $\phi(a)=\sigma\left(a^{\prime}\right) \in \mathscr{V}_{\phi} \cap \mathscr{V}_{\sigma}$. We complete the proof of Proposition 5.1 by summing the upper bound estimates over all choices of $d$ and all selections of the $u_{i}$ and $v_{i}$.

We carry out this plan in four stages, each of which is treated in more detail below:

- Finalize the notions of "small" and "tiny", and so also the choices of $d, e$, and $f$.
- Describe how to choose the $u_{i}$ and $v_{i}$ so that the intervals $\left[u_{i}, v_{i}\right]$ capture $P^{+}\left(p_{i}-1\right)$ and $P^{+}\left(q_{i}+1\right)$ for all $0 \leqslant i \leqslant k-1$.
- Check that the hypotheses of Lemma 4.1 are satisfied.
- Take the estimate of Lemma 4.1 and sum over $d$ and the choices of $u_{i}$ and $v_{i}$.

5A. "Small" and "tiny". Suppose we are given a solution $(a, a ') \in \mathscr{A}_{\phi} \times \mathscr{A}_{\sigma}$ to $\phi(a)=\sigma\left(a^{\prime}\right)$. Set $x_{j}=x_{j}(a ; x)$ and $y_{j}=x_{j}\left(a^{\prime} ; x\right)$, in the notation of Section 2B, so that (from the definition of $\left.\mathscr{A}_{f}\right)$ the sequences $\mathbf{x}=\left(x_{1}, \ldots, x_{L}\right)$ and $\mathbf{y}=\left(y_{1}, \ldots, y_{L}\right)$ belong to $\mathscr{S}_{L}(\xi)$.
Lemma 5.2. With $\left\{z_{j}\right\}_{j=1}^{L}$ denoting either of the sequences $\left\{x_{j}\right\}$ or $\left\{y_{j}\right\}$, we have
(i) $z_{j}<3 \varrho^{j}$ for $1 \leqslant j \leqslant L$,
(ii) $z_{L-j} \geqslant \frac{3}{100} \varrho^{-j} / \log _{2} x$ for $0 \leqslant j<L$.
(iii) $z_{j+2} \leqslant 0.9 z_{j}$ for $1 \leqslant j \leqslant L-2$.

Proof. Claim (i) is repeated verbatim from Lemma 2.8. By the same lemma, $z_{j} \leqslant 3 \varrho^{j-i} z_{i}$ for $1 \leqslant i<j \leqslant L$. This immediately implies (iii), since $\varrho^{2}<0.9$. Moreover, fixing $j=L$, condition (6) in the definition of $\mathscr{A}_{f}$ gives that

$$
z_{i} \geqslant \frac{1}{3} \varrho^{i-L} z_{L} \geqslant \frac{\log _{2} 3}{3} \varrho^{i-L} / \log _{2} x>\frac{3}{100} \varrho^{i-L} / \log _{2} x
$$

which is (ii) up to a change of variables.

Lemma 5.3. The minimal index $k_{0} \leqslant L$ for which

$$
\log _{2} P^{+}\left(p_{k_{0}}-1\right)<\left(\log _{2} x\right)^{1 / 2+\epsilon / 10}
$$

satisfies $k_{0} \sim(1 / 2-\epsilon / 10) L$ as $x \rightarrow \infty$.
Proof. Lemma 5.2(i) shows that the least $K$ with $\log _{2} p_{K}<\left(\log _{2} x\right)^{1 / 2+\epsilon / 10}$ satisfies $K \leqslant(1 / 2-\epsilon / 10+o(1)) L$, as $x \rightarrow \infty$. Since $\log _{2} P^{+}\left(p_{K}-1\right) \leqslant \log _{2} p_{K}$, this gives the asserted upper bound on $k_{0}$. The lower bound follows in a similar fashion from Lemma 5.2(ii) and Lemma 2.5.

Recall the definition of $\delta$ from (3-2), and put

$$
\eta:=10 L \delta, \quad \text { so that } \quad \eta \asymp\left(\log _{3} x\right)^{3 / 2}\left(\log _{2} x\right)^{-1 / 2} .
$$

We choose our "large"/"small" cutoff point $k$ by taking $k=k_{0}$ if $x_{k_{0}-1}-x_{k_{0}} \geqslant 20 \eta$, and taking $k=k_{0}-1$ otherwise. For future use, we note that with this choice of $k$,

$$
\begin{equation*}
x_{k-1}-x_{k} \geqslant 20 \eta . \tag{5-2}
\end{equation*}
$$

This inequality is immediate if $k=k_{0}$; in the opposite case, by Lemma 5.2(iii),

$$
\begin{aligned}
x_{k-1}-x_{k}=x_{k_{0}-2}-x_{k_{0}-1} & \geqslant x_{k_{0}-2}-x_{k_{0}}-20 \eta \\
& \geqslant 0.1 x_{k_{0}-2}-20 \eta \geqslant 0.1\left(\log _{2} x\right)^{-1 / 2+\epsilon / 10}-20 \eta>20 \eta .
\end{aligned}
$$

Note that with this choice of $k$, we have $\log _{2} p_{i}>\left(\log _{2} x\right)^{1 / 2+\epsilon / 10}$ for $0 \leqslant i \leqslant k-1$, and so condition (1) in the definition of $\mathscr{A}_{\phi}$ guarantees that each $p_{i}$ divides $a$ to the first power only, for $0 \leqslant i \leqslant k-1$. Moreover, from Lemmas 5.2(ii) and 5.3, we have $\log _{2} q_{i}>\left(\log _{2} x\right)^{1 / 2+\epsilon / 11}$ for $0 \leqslant i \leqslant k-1$. So each $q_{i}$ divides $a^{\prime}$ only to the first power, for $0 \leqslant i \leqslant k-1$. Now take

$$
f:=\phi\left(p_{k} p_{k+1} \cdots p_{L-1}\right), \quad d:= \begin{cases}\phi\left(p_{L} p_{L+1} \cdots\right) & \text { if } p_{L-1} \neq p_{L},  \tag{5-3}\\ \frac{p_{L}}{\phi\left(p_{L}\right)} \phi\left(p_{L} p_{L+1} \cdots\right) & \text { if } p_{L-1}=p_{L},\end{cases}
$$

and

$$
e:=\sigma\left(q_{k} q_{k+1} \cdots\right),
$$

and observe that (5-1) holds.
5B. Selection of the $\boldsymbol{u}_{\boldsymbol{j}}$ and $\boldsymbol{v}_{\boldsymbol{j}}$. Rather than choose the $u_{j}$ and $v_{j}$ directly, it is more convenient to work with the $\mu_{j}$ and $v_{j}$; then $u_{j}$ and $v_{j}$ are defined by $u_{j}:=\exp \left((\log x)^{\mu_{j}}\right)$ and $v_{j}:=\exp \left((\log x)^{v_{j}}\right)$. Put

$$
\begin{equation*}
\zeta_{0}:=1-\frac{\log _{3} x+\log 100}{\log _{2} x} \quad \text { and } \quad \zeta_{j}:=\zeta_{0}-j \eta \quad(j \geqslant 1), \tag{5-4}
\end{equation*}
$$

and note that with $\nu_{0}:=1$ and $\mu_{0}:=\zeta_{0}$, we have

$$
u_{0}=x^{1 /\left(100 \log _{2} x\right)}<x^{1 / \log _{2} x} \leqslant P^{+}\left(p_{0}-1\right), P^{+}\left(q_{0}+1\right) \leqslant x=v_{0},
$$

by condition (7) in the definitions of $\mathscr{A}_{\phi}$ and $\mathscr{A}_{\sigma}$. To choose the remaining $\mu_{j}$ and $v_{j}$, it is helpful to know that $p_{j}$ and $q_{j}$ are close together (renormalized on a double logarithmic scale) for $1 \leqslant j \leqslant k$. This is the substance of the following lemma.
Lemma 5.4. If $p_{j} \geqslant S$ and $q_{j} \geqslant S$, then $\left|x_{j}-y_{j}\right| \leqslant(2 j+1) \delta<\eta$. These hypotheses hold if $L-j \geqslant 2 C \log _{4} x+12$, and so in particular for $1 \leqslant j \leqslant k$.
Proof. Suppose for the sake of contradiction that $y_{j} \geqslant x_{j}+(2 j+1) \delta$; since the $p_{i}$ and $q_{i}$ are all $S$-normal, this would imply that

$$
(j+1)\left(y_{j}-x_{j}-\delta\right) \leqslant \frac{\Omega\left(\sigma\left(a^{\prime}\right), p_{j}, q_{j}\right)}{\log _{2} x}=\frac{\Omega\left(\phi(a), p_{j}, q_{j}\right)}{\log _{2} x} \leqslant j\left(y_{j}-x_{j}+\delta\right)
$$

which is false. We obtain a similar contradiction if we suppose $x_{j} \geqslant y_{j}+(2 j+1) \delta$. The second half of the lemma follows from Lemma 5.2 and a short calculation, together with the estimate $k \sim(1 / 2-\epsilon / 10) L$ of Lemma 5.3.

We choose the intervals $\left[\mu_{j}, v_{j}\right.$ ] for $1 \leqslant j \leqslant k-1$ successively, starting with $j=1$. (We select $v_{k}$ last, by a different method.) Say that the pair $\left\{x_{j}, x_{j+1}\right\}$ is well-separated if $x_{j}-x_{j+1} \geqslant 10 \eta$, and poorly separated otherwise.

In the well-separated case, among all $\zeta_{i}$ (with $i \geqslant 0$ ), choose $\zeta$ minimal and $\zeta^{\prime}$ maximal with

$$
\left.\left.\begin{array}{rl}
\zeta^{\prime} \log _{2} x \leqslant & \log _{2} \min \left\{P^{+}\left(p_{j}-1\right)\right.
\end{array}\right) P^{+}\left(q_{j}+1\right)\right\}, ~\left(\log _{2} \max \left\{P^{+}\left(p_{j}-1\right), P^{+}\left(q_{j}+1\right)\right\} \leqslant \zeta \log _{2} x, ~ \$\right.
$$

and put

$$
\mu_{j}:=\zeta, \quad v_{j}:=\zeta^{\prime}
$$

In the poorly separated case, $j<k-1$, by (5-2). We select [ $\left.\mu_{j}, v_{j}\right]=\left[\mu_{j+1}, v_{j+1}\right]$ by a similar recipe: Among all $\zeta_{i}$ (with $i \geqslant 0$ ), choose $\zeta$ minimal and $\zeta^{\prime}$ maximal with

$$
\begin{aligned}
& \zeta^{\prime} \log _{2} x \leqslant \log _{2} \min \left\{P^{+}\left(p_{j}-1\right), P^{+}\left(q_{j}+1\right), P^{+}\left(p_{j+1}-1\right), P^{+}\left(q_{j+1}+1\right)\right\} \\
& \leqslant \log _{2} \max \left\{P^{+}\left(p_{j}-1\right), P^{+}\left(q_{j}+1\right), P^{+}\left(p_{j+1}-1\right), P^{+}\left(q_{j+1}+1\right)\right\} \leqslant \zeta \log _{2} x
\end{aligned}
$$

and put

$$
v_{j}=v_{j+1}=\zeta, \quad \text { and } \quad \mu_{j}=\mu_{j+1}=\zeta^{\prime}
$$

To see that these choices are well-defined, note that by (7) in the definition of $\mathscr{A}_{f}$, we have $x_{j}, y_{j} \leqslant \zeta_{0}$, which implies that a suitable choice of $\zeta$ above exists in both cases. Also, for $1 \leqslant i \leqslant k$, we have $x_{i}, y_{i} \geqslant\left(\log _{2} x\right)^{-1 / 2+\epsilon / 11}$ (by Lemma 5.3 and 5.2(ii)). So by Lemma 2.5,

$$
\log _{2} \min \left\{P^{+}\left(p_{i}-1\right), P^{+}\left(q_{i}+1\right)\right\} / \log _{2} x \geqslant\left(\log _{2} x\right)^{-1 / 2+\epsilon / 12}
$$

say. Since neighboring $\zeta_{i}$ are spaced at a distance $\eta \asymp\left(\log _{2} x\right)^{-1 / 2}\left(\log _{3} x\right)^{3 / 2}$, a suitable choice of $\zeta^{\prime}$ also exists in both cases.

For our application of Lemma 4.1, it is expedient to keep track at each step of the length of the intervals $\left[\mu_{j}, v_{j}\right.$ ], as well as the distance between the left-endpoint of the last interval chosen and the right-endpoint of the succeeding interval (if any). In the well-separated case, Lemmas 5.4 and 2.5 show that

$$
v_{j} \leqslant \max \left\{x_{j}, y_{j}\right\}+\eta \leqslant x_{j}+2 \eta,
$$

while

$$
\begin{equation*}
\mu_{j} \geqslant \min \left\{x_{j}, y_{j}\right\}-\frac{\log _{3} x+\log 4}{\log _{2} x}-\eta \geqslant x_{j}-3 \eta, \tag{5-5}
\end{equation*}
$$

so that $\nu_{j}-\mu_{j} \leqslant 5 \eta$. Also, if a succeeding interval exists (so that $j+1 \leqslant k-1$ ), then $v_{j+1} \leqslant \max \left\{x_{j+1}, y_{j+1}\right\}+\eta \leqslant x_{j+1}+2 \eta$, and the separation between $\mu_{j}$ and $\nu_{j+1}$ satisfies the lower bound

$$
\begin{equation*}
\mu_{j}-v_{j+1} \geqslant x_{j}-x_{j+1}-5 \eta \geqslant 5 \eta . \tag{5-6}
\end{equation*}
$$

In the poorly separated case, we have

$$
v_{j} \leqslant \max \left\{x_{j}, y_{j}, x_{j+1}, y_{j+1}\right\}+\eta=\max \left\{x_{j}, y_{j}\right\}+\eta \leqslant x_{j}+2 \eta,
$$

as before, but the lower bound on $\mu_{j}$ takes a slightly different form;

$$
\begin{align*}
& \mu_{j} \geqslant \min \left\{x_{j}, y_{j}, x_{j+1}, y_{j+1}\right\}-\frac{\log _{3} x+\log 4}{\log _{2} x}-\eta \\
& \geqslant\left(x_{j+1}-\eta\right)-\frac{\log _{3} x+\log 4}{\log _{2} x}-\eta \geqslant x_{j+1}-3 \eta \geqslant x_{j}-13 \eta, \tag{5-7}
\end{align*}
$$

so that $v_{j}-\mu_{j} \leqslant 15 \eta$. In this case, since $v_{j}=v_{j+1}$ and $\mu_{j}=\mu_{j+1}$, the succeeding interval (if it exists) is [ $\mu_{j+2}, v_{j+2}$ ]. By Lemma 5.2(iii),

$$
x_{j}-x_{j+2} \geqslant 0.1 x_{j} \geqslant 0.1\left(\log _{2} x\right)^{-1 / 2+\epsilon / 10}>20 \eta,
$$

say. Thus, $v_{j+2} \leqslant \max \left\{x_{j+2}, y_{j+2}\right\}+\eta \leqslant x_{j+2}+2 \eta \leqslant x_{j}-18 \eta$, and so

$$
\begin{equation*}
\mu_{j+1}-v_{j+2}=\mu_{j}-v_{j+2} \geqslant\left(x_{j}-13 \eta\right)-\left(x_{j}-18 \eta\right) \geqslant 5 \eta . \tag{5-8}
\end{equation*}
$$

At this point we have selected intervals [ $\mu_{j}, \nu_{j}$ ], for all $0 \leqslant j \leqslant k-1$. We choose $\nu_{k}=\zeta$, where $\zeta$ is the minimal $\zeta_{i}$ satisfying $\zeta \geqslant x_{k}+\eta$. Note that

$$
\log _{2} S / \log _{2} x=36 \log _{3} x / \log _{2} x<\left(\log _{2} x\right)^{-1 / 2+\epsilon / 11} \leqslant x_{k}<\zeta=v_{k} \leqslant x_{k}+2 \eta
$$

Thus, $v_{k}>S$. From (5-5) and (5-7), $\mu_{k-1} \geqslant x_{k-1}-3 \eta$, so that also

$$
\begin{equation*}
\mu_{k-1}-v_{k} \geqslant x_{k-1}-x_{k}-5 \eta \geqslant 15 \eta, \tag{5-9}
\end{equation*}
$$

where the last estimate uses (5-2).

5C. Verification of hypotheses. We now check that Lemma 4.1 may be applied with $y=x$. By construction, $S \leqslant v_{k} \leqslant v_{k-1} \leqslant \cdots \leqslant v_{0}=x$, and $u_{i} \leqslant v_{i}$ for all $0 \leqslant i \leqslant k-1$. Moreover, if $\left[\mu_{j}, v_{j}\right] \neq\left[\mu_{j-1}, v_{j-1}\right]$ (where $2 \leqslant j \leqslant k-1$ ), then from (5-6) and (5-8), $\mu_{j-1}-v_{j} \geqslant 5 \eta=50 L \delta>2 \delta$, and from (5-9), $\mu_{k-1}-v_{k} \geqslant 15 \eta>2 \delta$. Thus, condition (a) of Lemma 4.1 is satisfied. It follows from our method of selecting the $\mu_{j}$ and $v_{j}$ that if $v_{j}=v_{j+1}$, then (again by (5-8)) $v_{j+2} \leqslant \mu_{j+1}-5 \eta<v_{j+1}=v_{j}$, which shows that condition (b) is also satisfied. Moreover, since $v_{k}>x_{k}$, we have $P^{+}(d) \leqslant p_{L} \leqslant p_{k}<v_{k}$. So we may focus our attention on hypotheses (i)-(v) of Lemma 4.1. We claim that these hypotheses are satisfied with our choices of $d, e$, and $f$ from Section 5A and with

$$
\begin{equation*}
l:=L-k \tag{5-10}
\end{equation*}
$$

Property (i) is contained in (2) from the definition of $\mathscr{A}_{f}$. By construction,

$$
u_{i} \leqslant P^{+}\left(p_{i}-1\right), P^{+}\left(q_{i}+1\right) \leqslant v_{i}
$$

for all $0 \leqslant i \leqslant k-1$, which is (ii). Since $v_{k} \geqslant S>\log y$, property (iii) holds by (1) in the definition of $\mathscr{A}_{f}$. The verification of (iv) is somewhat more intricate. Recalling that $v_{k}>x_{k}$, it is clear from (5-3) that

$$
P^{+}(f)<p_{k} \leqslant v_{k}
$$

To prove the same estimate for $P^{+}(e)$, we can assume $e \neq 1$. Let $r=P^{+}(e)$, and observe that $r \mid \sigma(R)$, for some prime power $R$ with $R \| q_{k} q_{k+1} \cdots$. If $R$ is a proper prime power, then from (1) in the definition of $\mathscr{A}_{f}$, we have

$$
r \leqslant \sigma(R) \leqslant 2 R \leqslant 2(\log x)^{2}<v_{k}
$$

So we can assume that $R$ is prime, and so $R \leqslant q_{k}$ and $r \leqslant P^{+}(R+1) \leqslant \max \{3, R\} \leqslant$ $q_{k}$. But by Lemma 5.4,

$$
\log _{2} q_{k} / \log _{2} x=y_{k} \leqslant x_{k}+(2 k+1) \delta<x_{k}+\eta \leqslant v_{k}
$$

Thus, $P^{+}(e)=r \leqslant v_{k}$. Hence, $P^{+}(e f) \leqslant v_{k}$. Turning to the second half of (iv), write $p_{k} \cdots p_{L-1}=A B$, where $A$ is squarefree, $B$ is squarefull and $\operatorname{gcd}(A, B)=1$. Recalling (2-2), we see that

$$
\Omega(\phi(A)) \leqslant 3 \Omega(A) \log _{2} v_{k} \leqslant 3 l \log _{2} v_{k}
$$

with $l$ as in (5-10). Let $B^{\prime}$ be the largest divisor of $a$ supported on the primes dividing $B$, so that $B^{\prime}$ is squarefull and $B \mid B^{\prime}$. By (1) in the definition of $\mathscr{A}_{f}$, we have $B^{\prime} \leqslant(\log x)^{2}$. If $B^{\prime} \leqslant \exp \left(\left(\log _{2} x\right)^{1 / 2}\right)$, then (estimating crudely)

$$
\Omega(\phi(B)) \leqslant \Omega\left(\phi\left(B^{\prime}\right)\right) \leqslant 2 \log \phi\left(B^{\prime}\right) \leqslant 2 \log B^{\prime} \leqslant 2\left(\log _{2} x\right)^{1 / 2}
$$

On the other hand, if $B^{\prime}>\exp \left(\left(\log _{2} x\right)^{1 / 2}\right)$, then by (4) in the definition of $\mathscr{A}_{f}$,

$$
\Omega(\phi(B)) \leqslant \Omega\left(\phi\left(B^{\prime}\right)\right) \leqslant 10 \log _{2} \phi\left(B^{\prime}\right) \leqslant 10 \log _{2} B^{\prime} \ll \log _{3} x .
$$

Since $\log _{2} v_{k}=v_{k} \log _{2} x>\eta \log _{2} x>\left(\log _{2} x\right)^{1 / 2}$, we have we have $\Omega(\phi(B)) \leqslant$ $2 \log _{2} v_{k}$ in either case. Hence,

$$
\Omega(f)=\Omega(\phi(A))+\Omega(\phi(B)) \leqslant(3 l+2) \log _{2} v_{k} \leqslant 4 l \log _{2} v_{k},
$$

which completes the proof of (iv). Finally, we prove (v): Suppose that $b \geqslant x^{1 / 3}$ is a divisor of $p_{0}-1$. Recalling again (2-2),

$$
P^{+}(b) \geqslant b^{1 / \Omega\left(p_{0}-1\right)} \geqslant b^{\frac{1}{3 \log _{2} x}} \geqslant x^{\frac{1}{9 \log _{2} x}}>x^{\frac{1}{100 \log _{2} x}} \geqslant v_{1} .
$$

Thus, setting $b$ to be the largest divisor of $p_{0}-1$ supported on the primes $\leqslant v_{1}$, we have $b<x^{1 / 3}$. From (3-3) and conditions (0) and (7) in the definition of $\mathscr{A}_{\phi}$,

$$
p_{0}=\frac{a}{p_{1} p_{2} p_{3} \ldots}>\frac{x / \log x}{x^{1 / 100} p_{1}}>x^{0.95},
$$

say. Thus, $\left(p_{0}-1\right) / b$ is a divisor of $p_{0}-1$ composed of primes greater than $v_{1}$ and of size at least $\left(p_{0}-1\right) x^{-1 / 3}>x^{9 / 10} x^{-1 / 3}>x^{1 / 2}$.

5D. Denouement. We are now in a position to establish Proposition 5.1 and so also Theorem 1.1. Suppose that $k$ and the $\mu_{i}$ and $\nu_{i}$ are fixed, as is $d$; this also fixes $l=L-k$. By Lemma 4.1, whose hypotheses were verified above, the number of values $\phi(a)$ coming from corresponding solutions to $\phi(a)=\sigma\left(a^{\prime}\right)$, with $\left(a, a^{\prime}\right) \in \mathscr{A}_{\phi} \times \mathscr{A}_{\sigma}$, is

$$
\begin{gather*}
\ll \frac{x}{d}\left(c \log _{2} x\right)^{6 k}(k+1)^{\Omega(d)}\left(\log v_{k}\right)^{8(k+l) \log (k+1)+1}(\log x)^{-2+\sum_{i=1}^{k-1} a_{i} v_{i}+E} \\
\leqslant \frac{x}{d} \exp \left(O\left(\left(\log _{3} x\right)^{2}\right)\right) L^{\Omega(d)}\left(\log v_{k}\right)^{L^{2}}(\log x)^{-2+\sum_{i=1}^{k-1} a_{i} x_{i}+E^{\prime}}, \tag{5-11}
\end{gather*}
$$

where

$$
E^{\prime}:=E+\sum_{i=1}^{k-1} a_{i}\left(v_{i}-x_{i}\right)=\delta \sum_{i=2}^{k}(i \log i+i)+2 \sum_{i=1}^{k-1}\left(v_{i}-\mu_{i}\right)+\sum_{i=1}^{k-1} a_{i}\left(v_{i}-x_{i}\right) .
$$

By our choice of $\nu_{i}$ and $\mu_{i}$ in Section 5B, we have $\nu_{i}-\mu_{i} \ll \eta$ and $\nu_{i}-x_{i} \ll \eta$. Hence,

$$
E^{\prime} \ll \delta L^{2} \log L+\eta\left(L+\sum_{i=1}^{k-1} a_{i}\right) \ll \delta L^{2} \log L+\eta L^{2} \log L \ll \delta L^{3} \log L
$$

In combination with (8) from the definition of $\mathscr{A}_{\phi}$, this shows that the exponent of $\log x$ on the right-hand side of $(5-11)$ is at most $-1-\omega+E^{\prime} \leqslant-1-\omega / 2$, and so

$$
(\log x)^{-2+\sum_{i=1}^{k-1} a_{i} x_{i}+E^{\prime}} \leqslant(\log x)^{-1} \exp \left(-\frac{1}{2}\left(\log _{2} x\right)^{1 / 2+\epsilon / 2}\right) .
$$

Moreover, by Lemma 5.3 and Lemma 5.2(i),

$$
\begin{equation*}
v_{k} \leqslant x_{k}+2 \eta \leqslant\left(\log _{2} x\right)^{-1 / 2+\epsilon / 9}+2 \eta \leqslant\left(\log _{2} x\right)^{-1 / 2+\epsilon / 5}, \tag{5-12}
\end{equation*}
$$

and hence

$$
\left(\log v_{k}\right)^{L^{2}}=\exp \left(L^{2}\left(\log _{2} x\right) v_{k}\right) \leqslant \exp \left(\left(\log _{2} x\right)^{1 / 2+\epsilon / 4}\right)
$$

Inserting all of this back into (5-11), we obtain an upper bound which is

$$
\begin{equation*}
\ll \frac{x}{\log x} \exp \left(-\frac{1}{3}\left(\log _{2} x\right)^{1 / 2+\epsilon / 2}\right) \frac{L^{\Omega(d)}}{d} . \tag{5-13}
\end{equation*}
$$

Now we sum over the parameters previously held fixed. We have $k<L$; also, for $i>0$, each $\mu_{i}$ and $v_{i}$ has the form $\zeta_{j}$ of (5-4). Thus, the number of possibilities for $k$ and the $\mu_{i}$ and $v_{i}$ is

$$
\begin{equation*}
\leqslant L\left(1+\left\lfloor\eta^{-1}\right\rfloor\right)^{2 L} \leqslant \exp \left(O\left(\left(\log _{3} x\right)^{2}\right)\right) \tag{5-14}
\end{equation*}
$$

Next, we prove that

$$
\begin{equation*}
\Omega(d) \ll\left(\log _{2} x\right)^{1 / 2} \tag{5-15}
\end{equation*}
$$

uniformly for the $d$ under consideration, so that

$$
\begin{equation*}
L^{\Omega(d)} \leqslant \exp \left(O\left(\left(\log _{2} x\right)^{1 / 2} \log _{4} x\right)\right) \tag{5-16}
\end{equation*}
$$

Put $m:=p_{L} p_{L+1} \cdots$. Suppose first that $p_{L} \neq p_{L-1}$, so that $m$ is a unitary divisor of $a$ and $d=\phi(m)$. If $m \leqslant \exp \left(\left(\log _{2} x\right)^{1 / 2}\right)$, then (5-15) follows from the crude bound $\Omega(d) \ll \log d$. On the other hand, if $m>\exp \left(\left(\log _{2} x\right)^{1 / 2}\right)$, then from (4) in the definition of $\mathscr{A}_{\phi}$, we have $\Omega(d)=\Omega(\phi(m)) \ll \log _{2} m$. But by (3) in the definition of $\mathscr{A}_{\phi}$ and Lemma 5.2(i),

$$
\begin{aligned}
\log _{2} m & \leqslant \log _{2} p_{L}^{10 \log _{2} x} \ll \log _{3} x+\log _{2} p_{L} \ll \log _{3} x+\varrho^{L} \log _{2} x \\
& \ll \log _{3} x+\varrho^{-2 \sqrt{\log _{3} x}} \varrho^{L_{0}} \log _{2} x \ll \varrho^{-2 \sqrt{\log _{3} x}} \log _{3} x \ll \exp \left(O\left(\sqrt{\log _{3} x}\right)\right),
\end{aligned}
$$

which again gives (5-15). Suppose now that $p_{L}=p_{L-1}$. In this case, let $m^{\prime}$ be the largest divisor of $a$ supported on the primes dividing $m$. Then $d \mid \phi\left(m^{\prime}\right)$, and so $\Omega(d) \leqslant \Omega\left(\phi\left(m^{\prime}\right)\right)$. Write $m^{\prime}=p_{L}^{j} m^{\prime \prime}$, where $j \geqslant 2$ and $p_{L} \nmid m^{\prime \prime}$; both $p_{L}^{j}$ and $m^{\prime \prime}$ are unitary divisors of $a$. We have $\Omega\left(\phi\left(m^{\prime \prime}\right)\right) \ll\left(\log _{2} x\right)^{1 / 2}$, by mimicking the argument used for $m$ in the case when $p_{L} \neq p_{L-1}$. Also, $\Omega\left(\phi\left(p_{L}^{j}\right)\right) \ll\left(\log _{2} x\right)^{1 / 2}$
except possibly if $p_{L}^{j}>\exp \left(\left(\log _{2} x\right)^{1 / 2}\right)$, in which case, invoking (1) and (4) in the definition of $\mathscr{A}_{\phi}$,

$$
\Omega\left(\phi\left(p_{L}^{j}\right)\right) \leqslant 10 \log _{2} \phi\left(p_{L}^{j}\right) \leqslant 10 \log _{2} p_{L}^{j} \leqslant 10 \log _{2}\left(\log ^{2} x\right) \ll \log _{3} x .
$$

So

$$
\Omega(d) \leqslant \Omega\left(\phi\left(m^{\prime}\right)\right)=\Omega\left(\phi\left(p_{L}^{j}\right)\right)+\Omega\left(\phi\left(m^{\prime \prime}\right)\right) \ll\left(\log _{2} x\right)^{1 / 2}
$$

confirming (5-15).
Referring back to (5-13), we see that it remains only to estimate the sum of $1 / d$. Since $P^{+}(d) \leqslant v_{k},(5-12)$ shows that every prime dividing $d$ belongs to the set $\mathscr{P}:=\left\{p: \log _{2} p \leqslant\left(\log _{2} x\right)^{1 / 2+\epsilon / 5}\right\}$. Thus,

$$
\begin{equation*}
\sum \frac{1}{d} \leqslant \prod_{p \in \mathscr{P}}\left(1+\frac{1}{p}+\frac{1}{p^{2}}+\cdots\right) \ll \exp \left(\left(\log _{2} x\right)^{1 / 2+\epsilon / 5}\right) \tag{5-17}
\end{equation*}
$$

Combining the estimates (5-13), (5-14), (5-16), and (5-17), we find that

$$
\#\left\{\phi(a): a \in \mathscr{A}_{\phi}, a^{\prime} \in \mathscr{A}_{\sigma}, \phi(a)=\sigma\left(a^{\prime}\right)\right\} \ll \frac{x}{\log x} \exp \left(-\frac{1}{4}\left(\log _{2} x\right)^{1 / 2+\epsilon / 2}\right)
$$

which completes the proof of Proposition 5.1 and of Theorem 1.1.

## References

[Canfield et al. 1983] E. R. Canfield, P. Erdős, and C. Pomerance, "On a problem of Oppenheim concerning "factorisatio numerorum"", J. Number Theory 17:1 (1983), 1-28. MR 85j:11012
[Erdős 1935] P. Erdős, "On the normal number of prime factors of $p-1$ and some related problems concerning Euler's $\phi$-function", Quart. Journ. of Math., Oxford Ser. 6 (1935), 205-213. Zbl 61.0129.05
[Erdös 1945] P. Erdös, "Some remarks on Euler's $\phi$-function and some related problems", Bull. Amer. Math. Soc. 51 (1945), 540-544. MR 7,49f Zbl 0061.08005
[Erdős 1959] P. Erdős, "Remarks on number theory, II: Some problems on the $\sigma$ function", Acta Arith. 5 (1959), 171-177. MR 21 \#6348 Zbl 0092.04601
[Erdős and Graham 1980] P. Erdős and R. L. Graham, Old and new problems and results in combinatorial number theory, Monographies de L'Enseignement Mathématique [Monographs of L'Enseignement Mathématique] 28, Université de Genève L'Enseignement Mathématique, Geneva, 1980. MR 82j:10001 Zbl 0434.10001
[Erdős and Hall 1973] P. Erdős and R. R. Hall, "On the values of Euler's $\varphi$-function", Acta Arith. 22 (1973), 201-206. MR 53 \#13143 Zbl 0252.10007
[Erdős and Hall 1976] P. Erdős and R. R. Hall, "Distinct values of Euler's $\phi$-function", Mathematika 23:1 (1976), 1-3. MR 54 \#2603 Zbl 0329.10036
[Ford 1998a] K. Ford, "The distribution of totients", Ramanujan Journal 2:1-2 (1998), 67-151. MR 99m:11106 Zbl 0914.11053
[Ford 1998b] K. Ford, "The distribution of totients", Electron. Res. Announc. Amer. Math. Soc. 4 (1998), 27-34. MR 99f:11125 Zbl 0888.11003
[Ford and Pollack 2011] K. Ford and P. Pollack, "On common values of $\phi(n)$ and $\sigma(m)$, I", Acta Math. Hungar. 133:3 (2011), 251-271. MR 2846095 Zbl 06006182
[Ford et al. 2010] K. Ford, F. Luca, and C. Pomerance, "Common values of the arithmetic functions $\phi$ and $\sigma "$, Bull. Lond. Math. Soc. 42:3 (2010), 478-488. MR 2011m:11191 Zbl 1205.11010
[Garaev 2011] M. Garaev, "On the number of common values of arithmetic functions $\varphi$ and $\sigma$ below $x$ ", Moscow J. Comb. Number Theory 1:3 (2011), 42-49. Zbl 06077912
[Halberstam and Richert 1974] H. Halberstam and H.-E. Richert, Sieve methods, London Math. Soc. Monographs 4, Academic Press, London, 1974. MR 54 \#12689 Zbl 0298.10026
[Hall and Tenenbaum 1988] R. R. Hall and G. Tenenbaum, Divisors, Cambridge Tracts in Mathematics 90, Cambridge University Press, 1988. MR 90a: 11107 Zbl 0653.10001
[Hardy and Ramanujan 1917] G. H. Hardy and S. Ramanujan, "The normal number of prime factors of a number n", Quart. J. Math 48 (1917), 76-92. Reprinted as pp. 262-275 in Collected papers of Srinivasa Ramanujan, edited by G. H. Hardy et al., Cambridge Univ. Press, 1927 (reprinted Chelsea, 2000). MR 2280878 Zbl 46.0262.03
[Maier and Pomerance 1988] H. Maier and C. Pomerance, "On the number of distinct values of Euler's $\phi$-function", Acta Arith. 49:3 (1988), 263-275. MR 89d:11083 Zbl 0638.10045
[Pillai 1929] S. S. Pillai, "On some functions connected with $\phi(n)$ ", Bull. Amer. Math. Soc. 35:6 (1929), 832-836. MR 1561819 JFM 55.0710.02
[Pomerance 1986] C. Pomerance, "On the distribution of the values of Euler's function", Acta Arith. 47:1 (1986), 63-70. MR 88b:11060 Zbl 0602.10035

Communicated by Andrew Granville
Received 2010-11-29 Revised 2011-11-30 Accepted 2012-01-30
ford@math.uiuc.edu Department of Mathematics, University of Illinois, 1409 West Green Street, Urbana, IL 61801, United States
pollack@uga.edu Department of Mathematics, University of Georgia, Boyd Graduate Studies Research Center, Athens, GA 30602, United States

# Algebra \& Number Theory 

msp.berkeley.edu/ant

## EDITORS

Managing Editor<br>Bjorn Poonen<br>Massachusetts Institute of Technology<br>Cambridge, USA

Editorial Board Chair
David Eisenbud
University of California
Berkeley, USA

## Board of Editors

| Georgia Benkart | University of Wisconsin, Madison, USA | Susan Montgomery | University of Southern California, USA |
| ---: | :--- | ---: | :--- |
| Dave Benson | University of Aberdeen, Scotland | Shigefumi Mori | RIMS, Kyoto University, Japan |
| Richard E. Borcherds | University of California, Berkeley, USA | Raman Parimala | Emory University, USA |
| John H. Coates | University of Cambridge, UK | Jonathan Pila | University of Oxford, UK |
| J-L. Colliot-Thélène | CNRS, Université Paris-Sud, France | Victor Reiner | University of Minnesota, USA |
| Brian D. Conrad | University of Michigan, USA | Karl Rubin | University of California, Irvine, USA |
| Hélène Esnault | Freie Universität Berlin, Germany | Peter Sarnak | Princeton University, USA |
| Hubert Flenner | Ruhr-Universität, Germany | Joseph H. Silverman | Brown University, USA |
| Edward Frenkel | University of California, Berkeley, USA | Michael Singer | North Carolina State University, USA |
| Andrew Granville | Université de Montréal, Canada | Vasudevan Srinivas | Tata Inst. of Fund. Research, India |
| Joseph Gubeladze | San Francisco State University, USA | J. Toby Stafford | University of Michigan, USA |
| Ehud Hrushovski | Hebrew University, Israel | Bernd Sturmfels | University of California, Berkeley, USA |
| Craig Huneke | University of Virginia, USA | Richard Taylor | Harvard University, USA |
| Mikhail Kapranov | Yale University, USA | Ravi Vakil | Stanford University, USA |
| Yujiro Kawamata | University of Tokyo, Japan | Michel van den Bergh | Hasselt University, Belgium |
| János Kollár | Princeton University, USA | Marie-France Vignéras | Université Paris VII, France |
| Yuri Manin | Northwestern University, USA | Kei-Ichi Watanabe | Nihon University, Japan |
| Barry Mazur | Harvard University, USA | Andrei Zelevinsky | Northeastern University, USA |
| Philippe Michel | École Polytechnique Fédérale de Lausanne | Efim Zelmanov | University of California, San Diego, USA |

PRODUCTION
production@msp.org
Silvio Levy, Scientific Editor

See inside back cover or www.jant.org for submission instructions.
The subscription price for 2012 is US $\$ 175 /$ year for the electronic version, and $\$ 275 /$ year ( $+\$ 40$ shipping outside the US) for print and electronic. Subscriptions, requests for back issues from the last three years and changes of subscribers address should be sent to Mathematical Sciences Publishers, Department of Mathematics, University of California, Berkeley, CA 94720-3840, USA.

Algebra \& Number Theory (ISSN 1937-0652) at Mathematical Sciences Publishers, Department of Mathematics, University of California, Berkeley, CA 94720-3840 is published continuously online. Periodical rate postage paid at Berkeley, CA 94704, and additional mailing offices.

ANT peer review and production are managed by EditFLOw ${ }^{\circledR}$ from Mathematical Sciences Publishers.
PUBLISHED BY

- mathematical sciences publishers
http://msp.org/
A NON-PROFIT CORPORATION
Typeset in $\mathrm{LAT}_{\mathrm{E}} \mathrm{X}$
Copyright ©2012 by Mathematical Sciences Publishers


## Algebra \& Number Theory

Volume 6 No. 82012
On the refined ramification filtrations in the equal characteristic case ..... 1579Liang Xiao
On common values of $\phi(n)$ and $\sigma(m)$, II ..... 1669
Kevin Ford and Paul Pollack
Galois representations associated with unitary groups over $\mathbb{Q}$ ..... 1697
Christopher Skinner
Abelian varieties and Weil representations ..... 1719
Sug Woo Shin
Small-dimensional projective representations of symmetric and alternating groups ..... 1773
Alexander S. Kleshchev and Pham Huu Tiep
Secant varieties of Segre-Veronese varieties ..... 1817
Claudiu Raicu


[^0]:    The first author was supported by NSF Grant DMS-0901339. The second author was supported by an NSF Postdoctoral Fellowship (award DMS-0802970). The research was conducted in part while the authors were visiting the Institute for Advanced Study, the first author supported by grants from the Ellentuck Fund and The Friends of the Institute For Advanced Study. Both authors thank the IAS for its hospitality and excellent working conditions. They also express their gratitude to the referee for a careful reading of the manuscript.

