# Computing Greatest Common Divisors and Factorizations in Quadratic Number Fields* 

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#### Abstract

In a quadratic number field $\mathbf{Q}(\sqrt{D}), D$ a squarefree integer, with class number 1 , any algebraic integer can be decomposed uniquely into primes, but for only 21 domains Euclidean algorithms are known. It was shown by Cohn [5] that for $D \leq-19$ even remainder sequences with possibly nondecreasing norms cannot determine the GCD of arbitrary inputs. We extend this result by showing that there does not even exist an input in these domains for which the GCD computation becomes possible by allowing nondecreasing norms or remainders whose norms are not as small as possible. We then provide two algorithms for computing the GCD of algebraic integers in quadratic number fields $\mathbf{Q}(\sqrt{D})$. The first applies only to complex quadratic number fields with class number 1 , and is based on a short vector construction in a lattice. Its complexity is $O\left(S^{3}\right)$, where $S$ is the number of bits needed to encode the input. The second algorithm allows us to compute GCD's of algebraic integers in arbitrary number fields (ideal GCD's if the class number is $>1$ ). It requires only $O\left(S^{2}\right)$ binary steps for fixed $D$, but works poorly if $D$ is large. Finally, we prove that in any domain, the computation of the prime factorization of an algebraic integer can be reduced in polynomial time to the problem of factoring its norm into rational primes. Our reduction is based on a constructive version of a theorem by A. Thue.


1. Introduction. The fact that the set of complex integers

$$
\{x+\sqrt{-1} y \mid x, y \in \mathbf{Z}\}
$$

forms a unique factorization domain (UFD) follows from the ability to perform divisions with remainder in this domain. Already C. F. Gauss generalized the complex integers to the domain of algebraic integers

$$
O_{d}=\left\{\left.\frac{x+\sqrt{d} y}{2} \right\rvert\, x, y \in \mathbf{Z}, x \equiv d y \bmod 2\right\}, \quad d= \begin{cases}4 D & \text { if } D \equiv 2,3 \bmod 4 \\ D & \text { if } D \equiv 1 \bmod 4\end{cases}
$$

of a quadratic number field $\mathbf{Q}(\sqrt{D}), D$ a squarefree integer, and proved that a Euclidean division (and thus unique factorization) was possible for $d=-11,-8,-7$, $-4,-3,5,8$ and several more positive discriminants $d$. In general, an abstract integral domain $R$ is Euclidean with respect to a degree function $N$ from the nonzero

[^0]elements of $R$ into the nonnegative integers if for any two elements $\alpha, \beta \in R, \beta \neq 0$, either $\beta$ divides $\alpha$, or there exists a (Euclidean) quotient $\gamma \in R$ and a (Euclidean) remainder $\rho \in R$ such that $\alpha=\beta \gamma+\rho$ and $N(\rho)<N(\beta)$. Once such a Euclidean division is constructible, the GCD of any two elements in $R$ can be determined by repeated division. The Euclidean algorithm (EA) consists of computing for any two elements $\rho_{0}, \rho_{1} \in R$ a sequence of Euclidean divisions,
$$
\rho_{i}=\rho_{i-2}-\gamma_{i-1} \rho_{i-1}, \quad i \geq 2
$$
such that $N\left(\rho_{i}\right)<N\left(\rho_{i-1}\right)$ or $\rho_{i}=0$, in which case $\operatorname{GCD}\left(\rho_{0}, \rho_{1}\right)=\rho_{i-1}$. In the cases of quadratic number fields mentioned before, the norm serves as a degree function for the Euclidean algorithm. In this paper we investigate the sequential complexity for GCD computations and prime factorizations in any quadratic number field, including non-Euclidean ones.

We measure the input size in terms of the rational and irrational parts of the inputs. Let

$$
R \xi=\frac{x}{2}, \quad I \xi=\frac{y \sqrt{|d|}}{2} \quad \text { for } \xi=\frac{x+y \sqrt{d}}{2} \in O_{d}
$$

and let

$$
\operatorname{size}(\xi)=\log (|R \xi|)+\log (|I \xi|)
$$

which is the number of bits necessary to write down $\xi$. In our notation, the norm of $\xi, \mathbf{N} \xi$, can be represented as follows:

$$
\mathbf{N} \xi=\xi \bar{\xi}=\frac{x+y \sqrt{d}}{2} \frac{x-y \sqrt{d}}{2}=\frac{x^{2}-d y^{2}}{4}=\left\{\begin{array}{ll}
(R \xi)^{2}-(I \xi)^{2} & \text { for } d>0 \\
(R \xi)^{2}+(I \xi)^{2} & \text { for } d<0
\end{array}\right\} \in \mathbf{Z}
$$

It is easy to show that $O_{d}$ is a Euclidean domain with respect to the norm for $d=-11,-8,-7,-4,-3$, but for no other $d<0$, and for $d=5,8,12$ and 13 . There are exactly twelve more discriminants $d>0$ for which $O_{d}$ is Euclidean with respect to the norm, namely $17,21,24,28,29,33,37,41,44,57,73$ and 76 (cf. Chatland and Davenport [4] and Barnes and Swinnerton-Dyer [1]).
$O_{d}$ is a UFD if and only if its class number $h(d)$ is one. There are exactly four more imaginary quadratic number fields with $h(d)=1$, namely $d=-19,-43,-67$, and -163 . That this list is exhaustive was finally established by Stark [29]. It is conjectured that infinitely many real quadratic number fields have class number 1. The list of those $d>0$ for which $O_{d}$ is not Euclidean with respect to the norm, but with $h(d)=1$, begins with $53,56,61,69,77,88,89,92,93, \ldots$.

By definition, the Euclidean algorithm (EA) with respect to the norm applies to $\rho_{0}, \rho_{1}$ exactly if there exists a norm-decreasing remainder sequence $\left\langle\rho_{1}, \ldots, \rho_{n}\right\rangle$ with

$$
\begin{align*}
\rho_{2}=\rho_{0}-\gamma_{1} \rho_{1}, \rho_{3}=\rho_{1}-\gamma_{2} \rho_{2}, \ldots, \rho_{n}=\rho_{n-2}-\gamma_{n-1} \rho_{n-1} &  \tag{1.1}\\
0 & =\rho_{n-1}-\gamma_{n} \rho_{n}
\end{align*}
$$

where $\gamma_{1}, \ldots, \gamma_{n} \in O_{d}$. The question arises, whether one can introduce another degree function with respect to which $O_{d}$ with $h(d)=1$ becomes Euclidean. Even more generally, one may drop the condition that the remainder sequence be decreasing with respect to any degree function and ask if there exists any sequence of the form (1.1), norm-decreasing or not. Before discussing this question, we introduce some more notation.

Generalizing (1.1), we consider a (finite or infinite) sequence $\Sigma=\left\langle\rho_{0}, \rho_{1}, \rho_{2}, \ldots\right\rangle$, where all $\rho_{i} \in \mathbf{Q}(\sqrt{d})$ (not necessarily $\in O_{d}$ ). We define $l=l(\Sigma)$ to be $\infty$, if $\Sigma$ is infinite, or $n+1$, if $\Sigma=\left\langle\rho_{0}, \ldots, \rho_{n+1}\right\rangle$. Following Cooke [6], we call $\Sigma$ a division chain, if
(i) for all $i, 1 \leq i<l(\Sigma)$, there exists $\gamma_{i} \in O_{d}$ such that $\rho_{i+1}=\rho_{i-1}-\gamma_{i} \rho_{i}$;
(ii) $\rho_{i} \neq 0$ for all $i<l$; if $l<\infty$, then $\rho_{l}=0$.

The Euclidean algorithm in the most general sense consists of computing a division chain $\left\langle\rho_{0}, \rho_{1}, \ldots, \rho_{n+1}\right\rangle$ and returning $\rho_{n}$, and every rule for specifying the choice of the remainders defines a version of (EA). If the norms of all remainders are minimized, that is, if for all $i<l, \gamma_{i}$ is chosen such that $\mathbf{N}\left(\rho_{i+1}\right)=$ $\mathbf{N}\left(\rho_{i-1}-\gamma_{i} \rho_{i}\right)$ has the smallest possible value, then we call $\Sigma$ a minimal remainder division chain, the corresponding instance of (EA) a minimal remainder version of (EA). Notice that we thus relax (1.1) in that we allow remainders with norm larger than the corresponding divisors. We also apply the attribute 'minimal remainder' to the individual divisions; thus the terms 'minimal remainder quotient,' 'minimal remainder,' 'minimal remainder division' have the obvious meaning. In this paper, the elements $\gamma_{i}$ will always have the meaning specified in (i), whenever a division chain $\Sigma$ is considered; for sequences $\Sigma^{\prime}, \Sigma^{\prime \prime}$, etc. we use $\gamma_{i}^{\prime}, \gamma_{i}^{\prime \prime}$ etc., without defining these numbers.

We now turn to the question raised above. First, consider the quadratic domains $O_{d}$ with $d \leq-19$. It is easy to show that the domains $O_{d}, d=-19,-43,-67$, and -163 are not Euclidean with respect to any degree function (Samuel [24]). Also, there exist $\rho_{0}, \rho_{1}$ such that no finite division chain beginning with $\rho_{0}, \rho_{1}$ exists (Cohn [5]). Thus there does not exist any version of (EA) which can compute $\operatorname{GCD}\left(\rho_{0}, \rho_{1}\right)$.

Our first group of results extends these facts. We show in Section 2 that if $\Sigma=$ $\left\langle\rho_{0}, \rho_{1}, \ldots\right\rangle$ is an arbitrary division chain and $\Sigma^{\prime}=\left\langle\rho_{0}, \rho_{1}, \rho_{2}^{\prime}, \rho_{3}^{\prime}, \ldots\right\rangle$ a minimal remainder division chain, then $l(\Sigma) \geq l\left(\Sigma^{\prime}\right)$. Thus, nothing can ever be gained by choosing remainders whose norms are not as small as possible; neither can this reduce the number of divisions needed, nor does there exist even a single input $\left(\rho_{0}, \rho_{1}\right)$ for which the computation of $\operatorname{GCD}\left(\rho_{0}, \rho_{1}\right)$ succeeds by choosing a nonminimal remainder at some stage, but does not succeed otherwise.

It has previously been shown for certain Euclidean domains that choosing minimal remainders with respect to the standard degree function minimizes the number of divisions. This was done by Lazard for $\mathbf{Z}$ and for $K[x], K$ a field ([17]), and also for $O_{-3}$ (private communication), and by Rolletschek [22] for $O_{-4}$. Here the standard degree function is the absolute value for the domain $\mathbf{Z}$, the degree in the usual sense for polynomials and the norm for imaginary quadratic number fields. For $d=-4$, that is the Gaussian integers, Caviness and Collins [3] have adopted Lehmer's idea for integer GCD (cf. Knuth [14, §4.5.2]), whereas Rolletschek [21], [22] established the equivalent of Lamés [15] bound on the maximum number of possible divisions necessary.

Our second theorem from Section 2 says that in $O_{d}, d=-19,-43,-67$, and $-163, \operatorname{GCD}\left(\rho_{0}, \rho_{1}\right)$ can be computed by some version of (EA) only if a normdecreasing sequence of remainders can be achieved. From this, the results of Samuel and Cohn follow as an easy corollary.

In the case $d>0$, the situation is different. Under a generalized Riemann hypothesis, Weinberger [32] shows that every unique factorization domain $O_{d}$, $d>0$, is Euclidean with respect to some degree function, and in Cooke and Weinberger [7] it is shown that a constant bound for the number of divisions can be achieved, namely 5 . In fact, these results are shown for the rings of algebraic integers in arbitrary algebraic number fields, provided there are infinitely many units. It is not shown, however, how one can efficiently construct these division chains and thus compute GCD's.

There remains the need for efficient algorithms for computing GCD's in those quadratic domains which are not Euclidean with respect to the norm, both real and imaginary. We will describe two such algorithms in this paper. The first, to be presented in Section 3, only applies to the imaginary quadratic domains $O_{d}$, $d=-19,-43,-67$ or -163 . It is based on a short vector construction in a lattice. The number of binary steps needed to compute $\operatorname{GCD}(\xi, \eta)$ is

$$
O\left(S^{3}\right), \quad S=\operatorname{size} \xi+\operatorname{size} \eta
$$

The polynomial bound would still remain valid for variable $d<0$. In other words, a polynomial upper bound for the complexity can be proved without using the fact shown in [29] that only 4 non-Euclidean imaginary quadratic domains $O_{d}$ satisfy $h(d)=1$.

The second algorithm, given in Section 4, is much more general, because it applies to all quadratic domains $O_{d}$, both real and imaginary, including those with $h(d)>1$. In the latter case, ideal GCD's are computed in a certain sense, which we will specify later. The algorithm has quadratic complexity for any fixed $d$. However, it requires some preparatory work (independent of $\xi, \eta$ ), which is quite costly if $d$ is large. It is not clear to us how to accomplish polynomial running time if $d$ is not fixed. Concluding the discussion of GCD-algorithms, we wish to point out that an asymptotically fast algorithm equivalent to Schönhage's [27] integer-half-GCDalgorithm is still not described at this time.

We next turn to factorization into primes. One realizes easily that an algorithm for factoring in $O_{d}$ can be devised which requires only an additional polynomial cost beyond the factorization of certain rational integers $x$. The reason is that a rational prime $p$ either remains irreducible in $O_{d}$, or factors into $p=\pi_{1} \pi_{2}$ with $\mathbf{N} \pi_{1}=\mathbf{N} \pi_{2}=p$. In this paper, we establish that from a factorization of $\mathbf{N} \xi$ into primes $p_{1} \cdots p_{k}$ we can construct for fixed $d$ in deterministic polynomial time the factorization $\xi=\pi_{1} \cdots \pi_{n}$. We solve this problem by using a constructive version of a theorem by A . Thue for solving the diophantine equation $x^{2}-d y^{2}=z p$ in $x, y$ and small $z$. This approach follows Shanks [28, Section 71], although our $z$ is about the square root of the one obtained there. The main feature of this approach is that it does not require the computation of greatest common divisors in $O_{d}$, which is the basis for another standard procedure to solve this problem. That our reduction is deterministic follows from a result by Schoof [26] and the fact that $d$ is fixed.
2. Properties of Quadratic Fields. The point of this section is to show that in imaginary quadratic domains $O_{d}, d \leq-19$, one cannot speed up (EA) or increase the set of inputs for which (EA) works by allowing remainder sequences which are not norm-decreasing or remainders of nonminimal norm. First we recall
some well-known facts about quadratic fields, the proofs of which can be found, e.g., in Hasse's text [9, §16]. Let

$$
\omega_{d}= \begin{cases}\frac{1+\sqrt{d}}{2} & \text { for } d \equiv 1 \bmod 4 \\ \frac{\sqrt{d}}{2} & \text { for } d \equiv 0 \bmod 4\end{cases}
$$

Then the set $\left\{1, \omega_{d}\right\}$ forms an integral basis for $O_{d}$, that is, every element of $O_{d}$ can be represented in the form $a \cdot 1+b \cdot \omega_{d}$, with $a, b \in \mathbf{Z}$.

A unit $\varepsilon \in O_{d}$ is an element such that $1 / \varepsilon \in O_{d}$. A necessary and sufficient condition for $\varepsilon$ to be a unit is that $|\mathbf{N} \varepsilon|=1$. For $d<0$, the multiplicative group of units is generated by $\{\sqrt{-1}\}$ for $d=-4,\{(1+\sqrt{-3}) / 2\}$ for $d=-3$, and $\{-1\}$ in all other cases. For $d>0$, there always exists a fundamental unit $\varepsilon_{1}, R \varepsilon_{1}>0, I \varepsilon_{1}>0$, such that the unit group is generated by $\left\{-1, \varepsilon_{1}\right\}$. Two elements $\xi_{1}, \xi_{2} \in O_{d}$ are associates, $\xi_{1} \sim \xi_{2}$, if there exists a unit $\varepsilon$ such that $\xi_{1}=\varepsilon \xi_{2}$. For $d>0$, any $\xi_{1}$ has, according to [9], an associate $\xi_{2}$ with

$$
\left|R \xi_{2}\right|,\left|I \xi_{2}\right|<\frac{\left|\mathbf{N} \xi_{1}\right|+\varepsilon_{1}}{2}
$$

However, we will need a better estimate in Section 4. By $\phi$ we denote the argument-function, which is defined for all complex numbers $\xi \neq 0$ by $\xi=|\xi|$. $(\cos (\phi(\xi))+i \sin (\phi(\xi)))$. We also need the analogous definition given in $[9$, p. 288] for elements of real quadratic fields: for $d>0, \phi(\xi)$ is defined by $\xi=$ $\operatorname{sign}(\xi) \sqrt{|\mathbf{N} \xi|} e^{\phi(\xi)}$; then $\phi(\bar{\xi})=-\phi(\xi)$. Dividing $\xi_{1}$ by a power $\varepsilon_{1}^{k}$ of $\varepsilon_{1}$ such that $\phi\left(\varepsilon_{1}^{k}\right)$ is as close to $\phi\left(\xi_{1}\right)$ as possible, we can also find an associate $\xi_{2}$ of $\xi_{1}$ such that $\left|\phi\left(\xi_{2}\right)\right| \leq \phi\left(\varepsilon_{1}\right) / 2$. By adding and subtracting the equalities

$$
\xi_{2}=\operatorname{sign}\left(\xi_{2}\right) \sqrt{\left|\mathbf{N} \xi_{2}\right|} e^{\phi\left(\xi_{2}\right)}, \quad \bar{\xi}_{2}=\operatorname{sign}\left(\bar{\xi}_{2}\right) \sqrt{\left|\mathbf{N} \xi_{2}\right|} e^{-\phi\left(\xi_{2}\right)}
$$

we get $2\left|R \xi_{2}\right|, 2\left|I \xi_{2}\right| \leq \sqrt{\left|\mathbf{N} \xi_{2}\right|} \cdot\left(e^{\phi\left(\xi_{2}\right)}+e^{-\phi\left(\xi_{2}\right)}\right)$, hence

$$
\left|R \xi_{2}\right|,\left|I \xi_{2}\right| \leq \sqrt{\left|\mathbf{N} \xi_{2}\right|} e^{\left|\phi\left(\xi_{2}\right)\right|} \leq \sqrt{\left|\mathbf{N} \xi_{2}\right|} \sqrt{\varepsilon_{1}}
$$

We now discuss how rational primes $p$ split in $O_{d}, h(d)=1$. If $p \mid d$, then $p \sim \pi^{2}$ for some prime $\pi \in O_{d}$. If $p \geq 3$ and the Legendre symbol $(d / p)=(D / p)=+1$, then there exists $l \in \mathbf{Z}$ such that $l^{2} \equiv D \bmod p$. Therefore, $p \mid(l+\sqrt{D})(l-\sqrt{D})$ and thus $p \sim \pi \bar{\pi}$ with $\pi=\operatorname{GCD}(p, l+\sqrt{D})$. If $p=2$ and $d=D \equiv 1 \bmod 8$, then

$$
2 \left\lvert\,\left(1-\frac{1+\sqrt{D}}{2}\right)\left(1-\frac{1-\sqrt{D}}{2}\right)\right.
$$

and thus $2 \sim \pi \bar{\pi}$ with $\pi=\operatorname{GCD}(2,1-(1+\sqrt{D}) / 2)$. In all other cases, $p$ is a prime in $O_{d}$.

We now come to the main results of this section, which apply to the domains $O_{d}, d=-19,-43,-67,-163$. Although one would usually expect that minimal remainder versions require the smallest number of divisions among all versions of (EA), one might suspect that there are exceptional instances, where other versions terminate faster. The following theorem shows, however, that this is not the case in those imaginary quadratic fields under consideration.

THEOREM 2.1. Let $\mathbf{Q}(\sqrt{d})$ be an imaginary quadratic number field whose corresponding number ring $O_{d}$ is a UFD without a Euclidean algorithm with respect to the norm, that is, $d$ has one of the values $-19,-43,-67,-163$. Let $\Sigma=\left\langle\rho_{0}, \rho_{1}, \rho_{2}, \ldots\right\rangle$ be an arbitrary division chain, and let $\Sigma^{\prime}=\left\langle\rho_{0}, \rho_{1}, \rho_{2}^{\prime}, \rho_{3}^{\prime}, \ldots\right\rangle$ be a minimal remainder division chain beginning with the same two elements $\rho_{0}, \rho_{1}$. Then $l(\Sigma) \geq l\left(\Sigma^{\prime}\right)$.

Proof. We apply a technique developed in Lazard [17] and in Rolletschek [22]. In what follows, we call a division chain $\Sigma=\left\langle\rho_{0}, \rho_{1} \ldots\right\rangle$ a counterexample, if there exists a sequence $\Sigma^{\prime}=\left\langle\rho_{0}, \rho_{1}, \rho_{2}^{\prime}, \ldots\right\rangle$ such that the assertion of the theorem does not hold. Since the theorem is trivially true for $l(\Sigma)=\infty$, we may apply induction on $l(\Sigma)$. Thus assume that $\Sigma=\left\langle\rho_{0}, \rho_{1}, \ldots, \rho_{n+1}\right\rangle$ is a counterexample, but that the theorem is true for every shorter sequence in place of $\Sigma$. Let $\Sigma^{\prime}$ be a minimal remainder division chain for $\rho_{0}, \rho_{1}$ with $l(\Sigma)<l\left(\Sigma^{\prime}\right)$. Then $n \geq 2$ and $\rho_{1} \neq \rho_{2}^{\prime}$, otherwise $\left\langle\rho_{1}, \rho_{2}, \ldots, \rho_{n+1}\right\rangle$ would be another counterexample, contrary to the induction hypothesis. Without loss of generality we may make the following three assumptions:
(i) All remainders in $\Sigma$, except possibly $\rho_{2}$, are minimal. For otherwise we could have considered a minimal remainder division chain $\left\langle\rho_{1}, \rho_{2}, \rho_{3}^{\prime \prime}, \ldots\right\rangle$ whose length $k$, by induction hypothesis, would have to be no more than $n$; but then $\left\langle\rho_{0}, \rho_{1}, \rho_{2}, \rho_{3}^{\prime \prime}, \ldots, \rho_{k+1}^{\prime \prime}\right\rangle$ would be another counterexample which would also satisfy (i).
(ii) $\gamma_{1}^{\prime}=0$; otherwise we could replace $\rho_{0}$ by $\rho_{0}-\gamma_{1}^{\prime} \rho_{1}, \gamma_{1}^{\prime}$ by 0 , to get a counterexample satisfying (ii).
(iii) $\rho_{1}=1$; otherwise we could divide $\rho_{0}, \ldots, \rho_{n+1}$ by $\rho_{1}$ to obtain a counterexample satisfying (iii). It is here that the consideration of nonintegral values of $\rho_{0}, \ldots$ makes the proof more convenient.

It follows from assumption (ii) that 0 is one of the algebraic integers in $O_{d}$ closest to $\rho_{0}=\rho_{2}$. Hence, $\rho_{0}$ must lie in a region $R_{1}=\{\alpha|\alpha \in \mathbf{Q}(\sqrt{d}),|\alpha| \leq|\alpha-\gamma|$ for all $\left.\gamma \in O_{d}\right\}$, which is shown in Figure 1. Since $\rho_{2} \neq \rho_{2}^{\prime}$, we have $\gamma_{1} \neq 0$. We consider the various possible values of $\gamma_{1}$. In several cases we have to consider the value of $\rho_{1} / \rho_{2}=1 / \rho_{2}$, and we denote this value by $\delta$.
(a) $\gamma_{1}=1$. Then $\rho_{2}$ lies in the region $R_{2}$ which is constructed by shifting $R_{1}$ to the left by $1 ; R_{2}$ is bounded by the straight lines $-3 / 2+y i$ and $-1 / 2+y i$ ( $y \in \mathbf{R}$ ), and four additional straight lines. Hence, $\delta$ certainly lies within the region $R_{2}^{\prime}$ shown in Figure 2; here the circles $C_{1}, C_{2}$ are the sets of inverses of all points of the form $-1 / 2+y i$ and of the form $-3 / 2+y i$, respectively.

Now recall the definition of the element $\omega_{d}$ at the beginning of this section. We need the fact that $I \omega_{d}>2$. $\delta$ has a distance $\leq 1$ from -1 , but a distance $>1$ from all lattice points outside the real axis. Hence $\gamma_{2}$, which is a minimal remainder quotient of $\rho_{1}$ and $\rho_{2}$ by assumption (i), and which is therefore one of the elements of $O_{d}$ closest to $\delta$, can only have one of the values $0,-1$ or -2 . Correspondingly, we have to consider three subcases.
(a1) $\gamma_{2}=0$. Then the sequence $\Sigma$ starts with $\left\langle\rho_{0}, 1, \rho_{0}-1,1, \rho_{0}, \rho_{5}, \ldots\right\rangle$, since the minimal remainder of $\rho_{0}-1$ and 1 equals the minimal remainder of $\rho_{0}$ and 1. (More precisely, we may assume without loss of generality that $\rho_{4}=\rho_{0}$ by the same argument as in the justification of assumption (i), although there may be several minimal remainders of $\rho_{0}$ and 1.) We now consider the division chain


Figure 1


Figure 2
$\Sigma^{\prime \prime \prime}=\left\langle\rho_{0}, 1, \rho_{0}, \rho_{5}, \ldots, \rho_{n+1}\right\rangle . \Sigma^{\prime \prime \prime}$ is actually a minimal remainder division chain, though not necessarily identical with $\Sigma^{\prime}=\left\langle\rho_{0}, 1, \rho_{0}, \rho_{3}^{\prime}, \ldots\right\rangle$. We can apply the induction hypothesis to the sequences formed from $\Sigma^{\prime \prime \prime}$ and $\Sigma^{\prime}$ by omitting their first elements; since $\Sigma^{\prime}$ is also a minimal remainder division chain, it follows that $l\left(\Sigma^{\prime}\right) \leq l\left(\Sigma^{\prime \prime \prime}\right)=n-1$, contradicting the assumption that $n+1=l(\Sigma)<l\left(\Sigma^{\prime}\right)$. This concludes the proof of the theorem for this case.
(a2) $\gamma_{2}=-1$. In this case, $\Sigma$ starts with $\left\langle\rho_{0}, 1, \rho_{0}-1, \rho_{0}, \rho_{4}, \ldots\right\rangle$. Again, we construct a division chain $\Sigma^{\prime \prime \prime}$ shorter than $\Sigma$ :

$$
\Sigma^{\prime \prime \prime}=\left\langle\rho_{0}, 1, \rho_{0},-\rho_{4}, \rho_{5}, \ldots,(-1)^{n} \rho_{n+1}\right\rangle
$$

this time using the fact that $1 \bmod \rho_{0}=-\left(\left(\rho_{0}-1\right) \bmod \rho_{0}\right)$. Then we can show as in (a1) that the given minimal remainder division chain $\Sigma^{\prime}$ satisfies $l\left(\Sigma^{\prime}\right) \leq l\left(\Sigma^{\prime \prime \prime}\right)<$ $l(\Sigma)$, a contradiction.
(a3) $\gamma_{2}=-2$. The sequence $\Sigma$ starts with

$$
\left\langle\rho_{0}, 1, \rho_{0}-1,2 \rho_{0}-1, \rho_{0}\left(1-2 \gamma_{3}\right)+\left(-1+\gamma_{3}\right), \rho_{5}, \rho_{6}, \ldots\right\rangle
$$

We put

$$
\Sigma^{\prime \prime \prime}=\left\langle\rho_{0}, 1, \rho_{0}, 1-2 \rho_{0}, \rho_{0}\left(-1+2 \gamma_{3}\right)+\left(1-\gamma_{3}\right),-\rho_{5},-\rho_{6}, \ldots\right\rangle
$$

choosing $\gamma_{3}^{\prime \prime \prime}=\gamma_{3}-1$. In this case, $l\left(\Sigma^{\prime}\right) \leq l\left(\Sigma^{\prime \prime \prime}\right)=l(\Sigma)$, a contradiction.
(b) $\gamma_{1}=-1$. This case parallels (a) completely, only some signs have to be changed.
(c) $\gamma_{1}=2$. As in case (a), we can determine the set of all possible values of $\delta$; it is the region $R_{3}^{\prime}$ shown in Figure 3.


Figure 3
It follows that $\gamma_{2}$ can only be 0 or -1 . The case $\gamma_{2}=0$ is treated as in case (a1). Assume $\gamma_{2}=-1 . \Sigma$ now has the form

$$
\left\langle\rho_{0}, 1, \rho_{0}-2, \rho_{0}-1, \rho_{0}\left(1-\gamma_{3}\right)+\left(-2+\gamma_{3}\right), \rho_{5}, \ldots,\right\rangle,
$$

and the sequence $\Sigma^{\prime \prime \prime}$ we construct has the form

$$
\left\langle\rho_{0}, 1, \rho_{0}, 1-\rho_{0}, \rho_{0}\left(-1+\gamma_{3}\right)+\left(2-\gamma_{3}\right),-\rho_{5}, \ldots,\right\rangle,
$$

where we choose $\gamma_{3}^{\prime \prime \prime}=-2+\gamma_{3}$. Then $l\left(\Sigma^{\prime \prime \prime}\right)=l(\Sigma)$, and the rest of the proof parallels previous cases.
(d) $\gamma_{1}=-2$. This case is analogous to (c).
(e) $\gamma_{1}$ is real, $\left|\gamma_{1}\right| \geq 3$. Then $\left|\rho_{2}\right|=\left|\rho_{0}-\gamma_{1}\right|>2$, hence $|\delta|<1 / 2$. Then $\gamma_{2}=0$, and the assertion of the theorem follows as in (a1).
(f) $\gamma_{1}$ is not real: $\gamma_{1}=a+b i$ with $b \neq 0$. We show in all cases $\left|I\left(\rho_{2}\right)\right|=$ $\left|I\left(\rho_{0}-\gamma_{1}\right)\right|>1$. Indeed, the minimal absolute value of $I\left(\rho_{0}-\gamma_{1}\right)$ occurs with $d=-19, \gamma_{1}=\omega_{d}$ and $\rho_{0}=\varsigma$, as pictured in Figure 1, and in this case $\left|I\left(\rho_{2}\right)\right|>1$.


Figure 4
Figure 4 shows the region $R_{4}^{\prime}$, which is the set of all inverses of complex numbers $\alpha$ with $|I(\alpha)|>1$, and which contains the set of all possible values of $\delta$. Again it follows that the element $\gamma_{2} \in O_{d}$ closest to $\delta$ can only be 0 , so $\Sigma=\left\langle\rho_{0}, 1, \rho_{0}\right.$ $\left.\gamma_{1}, \rho_{4}, \ldots\right\rangle$. But now, $\Sigma^{\prime \prime \prime}=\left\langle\rho_{0}, 1, \rho_{4}, \ldots\right\rangle$ is a shorter division chain than $\Sigma$, leading to the same contradiction as in the previous cases.

The cases (a)-(f) are exhaustive, completing the proof.
The condition $d \leq-19$ in the previous theorem is used in the estimates in the cases (a) and (f) in the above proof. Remarkably, the statement of the theorem (without the assumption of $O_{d}$ being non-Euclidean) fails for $d=-11$, but it remains valid for all other Euclidean imaginary number fields, as was recently shown in [23].

In algorithm-theoretic terms, Theorem 2.1 can be formulated as follows:
Corollary 2.1. Let $d$ be as in Theorem 2.1. Then:
(i) All minimal remainder versions of (EA), applied to $\rho_{0}, \rho_{1}$, require the same number of divisions, and no other version requires fewer.
(ii) If no minimal remainder version of (EA) allows the computation of $\operatorname{GCD}\left(\rho_{0}, \rho_{1}\right)$, then no division chain terminates with $\operatorname{GCD}\left(\rho_{0}, \rho_{1}\right)$.

While Theorem 2.1 and its corollary deal with divisions where the remainder has nonminimal norm, the next theorem shows that nothing is gained by allowing divisions where the remainder has norm greater than the divisor, even if that remainder is minimal. In other words, if a minimal remainder version of (EA), applied to $\rho_{0}, \rho_{1}$, leads to a division where no remainder with norm smaller than the divisor exists, then every version of (EA) fails for this input.

THEOREM 2.2. For any $d \leq-19, h(d)=1$, and for all $\rho_{0}, \rho_{1} \in O_{d}$, there exists a sequence $\gamma_{1}, \ldots, \gamma_{n} \in O_{d}$ satisfying (1.1) only if a norm-decreasing sequence with this property exists, that is, if the common version of (EA) applies to $\rho_{0}, \rho_{1}$, where the remainder of each division has smaller norm than the divisor.

Proof. Consider a minimal remainder division chain $\Sigma=\left\langle\rho_{0}, \rho_{1}, \rho_{2}, \ldots\right\rangle$, and assume that $\left\langle\rho_{1}, \rho_{2}, \ldots\right\rangle$ is not norm-decreasing. We will show that there exists an infinite minimal remainder division chain $\Sigma^{\prime}$ starting with $\rho_{0}, \rho_{1}$. Similarly as in the proof of Theorem 2.1, a number of assumptions can be made without loss of generality:
(i) $\left|\mathbf{N} \rho_{2}\right| \geq\left|\mathbf{N} \rho_{1}\right|$;
(ii) $\rho_{1}=1$;
(iii) $\gamma_{1}=0$;
(iv) $R \rho_{0} \geq 0, I \rho_{0} \geq 0$.

The justification of (i) is immediate. For (ii) and (iii), see the analogous assumptions in the proof of Theorem 2.1. Finally, (iv) can be justified by symmetry; in the following proof only some signs and limits for $\phi\left(\rho_{0}\right)$ would have to be changed if $R \rho_{0}$ and/or $I \rho_{0}$ is $<0$. We note in passing that an assumption analogous to (iv) could also have been made in the proof of Theorem 2.1, but its justification would have required a few lines, and the rest of the proof would not have been simplified.

We have $\rho_{2}=\rho_{0}$ by (iii), $\mathbf{N}\left(\rho_{0}\right)=\mathbf{N}\left(\rho_{2}\right) \geq 1$ by (i) and (ii). By (iii), 0 is one of the algebraic integers in $O_{d}$ closest to $\rho_{0}=\rho_{2}$. Together, these facts imply that $\rho_{0}$ lies in the region shown in Figure 5.


Figure 5
It follows now that $\pi / 3 \leq \phi\left(\rho_{0}\right) \leq \pi / 2$, hence $\delta=1 / \rho_{0}$ satisfies $-\pi / 2 \leq \phi(\delta) \leq$ $-\pi / 3 ;$ also, $\mathbf{N}(\delta) \leq 1$. Then one of the algebraic integers $\gamma_{2} \in O_{d}$ closest to $\delta$ is 0 . Choosing $\gamma_{2}=0$, we get $\rho_{3}=1$. Hence, the following is a minimal remainder division chain:

$$
\Sigma^{\prime}=\left\langle\rho_{0}, 1, \rho_{0}, 1, \rho_{0}, \ldots\right\rangle
$$

This sequence is infinite, as desired. By Theorem 2.1, every division chain starting with $\rho_{0}, \rho_{1}$ is infinite, so no version of (EA) can terminate.

COROLLARY 2.2. The domain $O_{d}, d \leq-19, h(d)=1$, is not Euclidean for any choice of degree function.

Proof. Recall that for $d \leq-19, O_{d}$ is not Euclidean with respect to the norm. This means, by definition, that there exist elements $\rho_{0}, \rho_{1} \in O_{d}$ such that every
remainder of $\rho_{0}$ and $\rho_{1}$ has larger norm than $\rho_{1}$. Thus, if $\Sigma$ is a minimal remainder division chain beginning with $\rho_{0}, \rho_{1}$, then $\left\langle\rho_{1}, \rho_{2}, \ldots\right\rangle$ is not norm-decreasing, and by Theorem 2.2 no finite division chain beginning with $\rho_{0}, \rho_{1}$ exists. The assertion follows.
3. GCD Computation by Lattice Reduction. As we have seen in Section 2, not all quadratic fields with unique factorization allow a Euclidean algorithm, e.g., $\mathbf{Q}(\sqrt{-19})$ and $\mathbf{Q}(\sqrt{53})$. Therefore, a different GCD procedure is required for these domains. Our first algorithm for imaginary quadratic fields is based on computing short integral lattice vectors and is interesting for two reasons, even though the algorithm in Section 4 has lower asymptotic complexity. First, it does not require any field-dependent preconditioning, and second, its running time is polynomial, independently of the fact that $|d|$ is known to be bounded. The main idea of the algorithm is to solve $\lambda \xi+\mu \eta=0$ such that $\mathbf{N} \mu$ is small. The following lemma will be useful.

Lemma 3.1. Let $\xi, \eta \in O_{d}, d \leq-19, h(d)=1, \xi \eta \neq 0, \delta=\operatorname{GCD}(\xi, \eta)$, $\xi^{*}=\xi / \delta, \eta^{*}=\eta / \delta$. Assume that $\lambda \xi+\mu \eta=0, \lambda, \mu \in O_{d}$. Then $\eta^{*}\left|\lambda, \xi^{*}\right| \mu$. Furthermore, if $\lambda$ is not an associate of $\eta^{*}$, then $\mathbf{N} \mu \geq 4 \mathbf{N} \xi^{*}$.

Proof. Since $\lambda \xi^{*}=-\mu \eta^{*}$ and $\operatorname{GCD}\left(\xi^{*}, \eta^{*}\right)=1$, all prime factors of $\eta^{*}$ must occur in $\lambda$, and similarly for $\xi^{*}$ and $\mu$. Since 2 and 3 are primes in any $O_{d}$ in question, $\mathbf{N}\left(\mu / \xi^{*}\right) \geq 4$.

We consider the case $d \equiv 1 \bmod 4$ only, since only this one occurs. Let $\lambda=$ $\left(l_{1}-l_{2} / 2\right)+\sqrt{d} l_{2} / 2, \xi=\left(x_{1}-x_{2} / 2\right)+\sqrt{d} x_{2} / 2, \mu=\left(m_{1}-m_{2} / 2\right)+\sqrt{d} m_{2} / 2$, $\eta=\left(y_{1}-y_{2} / 2\right)+\sqrt{d} y_{2} / 2$ with $l_{1}, l_{2}, m_{1}, m_{2}, x_{1}, x_{2}, y_{1}, y_{2} \in \mathbf{Z}$. Then

$$
\begin{align*}
4 R(\lambda \xi+\mu \eta)= & \left(4 x_{1}-2 x_{2}\right) l_{1}+\left(x_{2}-2 x_{1}+d x_{2}\right) l_{2} \\
& +\left(4 y_{1}-2 y_{2}\right) m_{1}+\left(y_{2}-2 y_{1}+d y_{2}\right) m_{2}  \tag{3.1}\\
\frac{2 I(\lambda \xi+\mu \eta)}{\sqrt{d}}= & x_{2} l_{1}+\left(x_{1}-x_{2}\right) l_{2}+y_{2} m_{1}+\left(y_{1}-y_{2}\right) m_{2}
\end{align*}
$$

We want to find integers $l_{1}, l_{2}, m_{1}$ and $m_{2}$ such that the right sides of the equations (3.1) are 0 , meaning that for the corresponding $\lambda$ and $\mu, \lambda \xi+\mu \eta=0$. Simultaneously, we want to keep $\mathbf{N} \mu$ small because the smallest such $\mu$ is an associate of $\xi^{*}$ by the previous lemma. This leads to the problem of finding a short vector in an integer lattice. The next theorem shows that only associates of $\xi^{*}$ can correspond to short vectors in a particular lattice, and thus finding a short vector in that lattice with a reduction algorithm actually gives $\xi^{*}$ and hence $\delta=\xi / \xi^{*}$.

THEOREM 3.1. Let $O_{d}$ be a UFD with $d \leq-19$, and let $\xi=\left(x_{1}-x_{2} / 2\right)+$ $\sqrt{d} d x_{2} / 2$ and $\eta=\left(y_{1}-y_{2} / 2\right)+\sqrt{d} y_{2} / 2$ be two nonzero elements of $O_{d}$. Let $c$ and $\hat{d}$ be integers such that $c \geq \sqrt{12 \mathbf{N} \xi},|\sqrt{|d|}-\hat{d}| \leq 1 / 2$. Furthermore, let $L^{*}$ be the 4-dimensional integer lattice spanned by the columns of the matrix

$$
L=\left[\begin{array}{cccc}
c\left(4 x_{1}-2 x_{2}\right) & c\left(x_{2}-2 x_{1}+d x_{2}\right) & c\left(4 y_{1}-2 y_{2}\right) & c\left(y_{2}-2 y_{1}+d y_{2}\right) \\
c x_{2} & c\left(x_{1}-x_{2}\right) & c y_{2} & c\left(y_{1}-y_{2}\right) \\
0 & 0 & 2 & -1 \\
0 & 0 & 0 & \hat{d}
\end{array}\right]
$$

Then, if $\nu^{*}=L \times\left[l_{1}, l_{2}, m_{1}, m_{2}\right]^{T} \neq 0$ is a vector of shortest Euclidean length in the lattice $L^{*}$, then $\mu=\left(m_{1}-m_{2} / 2\right)+\sqrt{d} m_{2} / 2$ must be an associate of $\xi^{*}=$ $\xi / \operatorname{GCD}(\xi, \eta)$. Furthermore, if $\nu \in L^{*} \backslash\{0\}$ does not have shortest Euclidean length, then $\|\nu\|^{2}>\frac{12}{5}\left\|\nu^{*}\right\|^{2}$.

Proof. Since

$$
\operatorname{det}\left[\begin{array}{cc}
4 x_{1}-2 x_{2} & x_{2}-2 x_{1}+d x_{2} \\
x_{2} & x_{1}-x_{2}
\end{array}\right]=\left(2 x_{1}+x_{2}\right)^{2}-d x_{2}^{2}>0, \quad x_{1} x_{2} \neq 0,
$$

$L$ is of full rank. With $\lambda$ and $\mu$ as above, we get from (3.1) that

$$
L \times\left[\begin{array}{c}
l_{1} \\
l_{2} \\
m_{1} \\
m_{2}
\end{array}\right]=\left[\begin{array}{c}
4 c R(\lambda \xi+\mu \eta) \\
2 c I(\lambda \xi+\mu \eta) / \sqrt{d} \\
2 m_{1}-m_{2} \\
\hat{d} m_{2}
\end{array}\right]=\nu=\left[\begin{array}{c}
\nu_{1} \\
\nu_{2} \\
\nu_{3} \\
\nu_{4}
\end{array}\right] \in \mathbf{Z}^{4} .
$$

We now estimate the Euclidean length $\|\nu\|$ of $\nu \neq 0$. If $\lambda \xi+\mu \eta \neq 0$ then $\|\nu\|^{2} \geq$ $c^{2} \geq 12 \mathbf{N} \xi$. In case $\lambda \xi+\mu \eta=0$, we have $\nu_{1}=\nu_{2}=0$, and we get from $\left||d|-\hat{d}^{2}\right| \leq \hat{d}$, $\hat{d} /|d|<1 / 4, d \leq-19$, and $m_{2}^{2} \leq 4 \mathrm{~N} \mu /|d|$ that

$$
\begin{align*}
\left|\|\nu\|^{2}-4 \mathbf{N} \mu\right| & =\left|\left(2 m_{1}-m_{2}\right)^{2}+\hat{d}^{2} m_{2}^{2}-\left(\left(2 m_{1}-m_{2}\right)^{2}-d m_{2}^{2}\right)\right| \\
& =\left|\hat{d}^{2}-|d|\right| m_{2}^{2} \leq \hat{d} m_{2}^{2} \leq \frac{4 \hat{d} \mathbf{N} \mu}{|d|}<\mathbf{N} \mu \tag{3.2}
\end{align*}
$$

From (3.2) we conclude that

$$
\begin{equation*}
3 \mathbf{N} \mu<\|\nu\|^{2}<5 \mathbf{N} \mu \quad \text { for } \lambda \xi+\mu \eta=0 . \tag{3.3}
\end{equation*}
$$

Therefore, under the assumption that $\mu$ is an associate of $\xi^{*}$ and $\lambda \xi+\mu \eta=0$, the corresponding vector $\nu^{*}$ satisfies $\left\|\nu^{*}\right\|^{2}<5 \mathbf{N} \mu=5 \mathbf{N} \xi^{*}$. Otherwise, Lemma 3.1 states that $\mathbf{N} \mu \geq 4 \mathbf{N} \xi^{*}$, and thus we get from (3.3)

$$
\|\nu\|^{2} \geq \min (12 \mathbf{N} \xi, 3 \mathbf{N} \mu) \geq 12 \mathbf{N} \xi^{*}
$$

Therefore,

$$
\|\nu\|^{2} \geq 12 \mathbf{N} \xi^{*}>\frac{12}{5}\left\|\nu^{*}\right\|^{2}
$$

which proves the theorem.
The algorithm for computing the GCD of $\xi$ and $\eta$ is now easy. One computes a vector in the lattice $L^{*}$ whose length is within a factor $C^{3 / 2}$ of the length of the shortest vector, where $C$ is a constant with $C>4 / 3$ and $C^{3}<12 / 5$, e.g., $C=83 / 62$. There exist several versions of the basis reduction algorithm by A. K. Lenstra et al. [18], cf. [12]. Already the original algorithm [18] can find such a vector. Since the dimension is fixed, all these algorithms take $O\left(\operatorname{size}^{3} \xi \eta\right)$ binary steps to compute such a vector (slightly less, if fast multiplication is used). But by Theorem 3.1 such a vector must be a shortest vector whose entries determine an associate of $\xi^{*}$. We then obtain $\delta=\operatorname{GCD}(\xi, \eta)=\xi / \xi^{*}$.
4. Greatest Common Divisor Computation. We present a GCD procedure that works for all quadratic domains. It turns out that one can always divide an appropriate multiple $l \xi, l \in \mathbf{Z}$, of the dividend $\xi$ by $\eta$ and accomplish a remainder of norm smaller than $|\mathbf{N} \eta|$. Moreover, the size of $l$ only depends on $d$ and not on $\xi$ or $\eta$. This fact was shown for arbitrary number fields by Hurwitz, see [11, p. 237]. The following lemma provides the specific bound for $l$ for quadratic number fields $O_{d}$. It is not restricted to the case $h(d)=1$.

Lemma 4.1. Let $c \in \mathbf{R}, 1 / 2 \leq c \leq 1, O_{d}$ be the ring of integers in a quadratic number field. Then for all $\xi, \eta \in O_{d}$, there exists an $l$ such that

$$
l \in \mathbf{Z} \quad \text { with } 1 \leq l \leq\left\lfloor\frac{\sqrt{|d|}}{2 c}\right\rfloor
$$

and there exists a $\gamma \in O_{d}$ with

$$
\begin{equation*}
|R(l \xi / \eta-\gamma)| \leq \frac{1}{2}, \quad|I(l \xi / \eta-\gamma)|<c \tag{4.1}
\end{equation*}
$$

Furthermore, if $d>0$, then $|\mathbf{N}(l \xi-\gamma \eta)| \leq c^{2}|\mathbf{N} \eta|$; if $d<0$, then $|\mathbf{N}(l \xi-\gamma \eta)|<$ $\left(c^{2}+1 / 4\right)|\mathbf{N} \eta|$.

Proof. This follows from the theory of approximation of real by rational numbers and the theory of continued fractions. We apply Theorem 171 in Hardy and Wright [8]: if $p_{n} / q_{n}$ and $p_{n+1} / q_{n+1}$ are the $n$th and $(n+1)$ st continued fraction approximation of a real number $x$, then

$$
\left|x-\frac{p_{n}}{q_{n}}\right|<\frac{1}{q_{n} q_{n+1}}<\frac{1}{q_{n}^{2}}
$$

We apply this theorem to $x=I(\xi / \eta) / I\left(\omega_{d}\right)=I(\xi / \eta) /(\sqrt{|d|} / 2)$, and a continued fraction approximation $p_{n} / q_{n}$ of $x$ such that $q_{n} \leq \sqrt{|d|} / 2 c$ and either $p_{n} / q_{n}=x$ or $q_{n+1}>\sqrt{|d|} / 2 c$. Then

$$
\left|\frac{I(\xi / \eta)}{\sqrt{|d|} / 2}-\frac{p_{n}}{q_{n}}\right|<\frac{1}{q_{n} q_{n+1}}<\frac{c}{q_{n} \sqrt{|d|} / 2}
$$

hence

$$
\left|I\left(q_{n} \xi / \eta-p_{n} \frac{\sqrt{|d|}}{2}\right)\right|=\left|q_{n} I(\xi / \eta)-p_{n} \frac{\sqrt{|d|}}{2}\right|<c .
$$

Then (4.1) is satisfied for $l=q_{n}, \gamma=p_{n} \omega_{d}+\left\lfloor R\left(q_{n} \xi / \eta-p_{n} \omega_{d}+1 / 2\right)\right\rfloor$.
Now let $\tau=l \xi / \eta-\gamma$. For $d>0$,

$$
|\mathbf{N} \tau|=\left|R(\tau)^{2}-I(\tau)^{2}\right| \leq \max \left(R(\tau)^{2}, I(\tau)^{2}\right)=c^{2}
$$

For $d<0$,

$$
|\mathbf{N} \tau|=R(\tau)^{2}+I(\tau)^{2} \leq c^{2}+1 / 4
$$

Then the second statement of the lemma follows by multiplying these inequalities by $|\mathbf{N} \eta|$.

This lemma suggests the following algorithm for computing GCD's in arbitrary quadratic domains. We formulate it first for those $d$ with unique factorization. At the end of this section we will discuss how to adapt it to apply to the case $h(d)>1$.
Preconditioned GCD in $O_{d}$.
Input: $\xi, \eta \in O_{d}, O_{d}$ a UFD.
Preconditioning: Given are the prime factors of $l=2, \ldots, m=\lfloor\sqrt{|d|}\rfloor$ in $O_{d}$,

$$
l=\pi_{1}^{e_{l, 1}} \cdots \pi_{k}^{e_{l, k}}, \quad 2 \leq l \leq m, e_{l, j} \geq 0
$$

If $d>0$ we furthermore are given a fundamental unit $\varepsilon_{1}$.
Output: $\delta=\operatorname{GCD}(\xi, \eta)$.

Step 1: $\Pi \leftarrow \varnothing ; \rho_{0} \leftarrow \xi ; \rho_{1} \leftarrow \eta$.
FOR $i \leftarrow 1,2, \ldots$ WHILE $\rho_{i} \neq 0$ DO step 2.
Step 2: Here we carry out a "pseudo-remainder" step, that is, we compute $l_{i+1} \in \mathbf{Z}$, $1 \leq l_{i+1} \leq m$, and $\gamma_{i+1} \in O_{d}$ such that

$$
\left|\mathbf{N}\left(l_{i+1} \rho_{i-1}-\gamma_{i+1} \rho_{i}\right)\right|<\frac{1}{2}\left|\mathbf{N} \rho_{i}\right|
$$

Set $\left(x_{i}+y_{i} \sqrt{d}\right) / 2 \leftarrow \rho_{i-1} / \rho_{i}$ where $x_{i}, y_{i} \in \mathbf{Q}$.
Determine $l_{i+1}$ satisfying Lemma 3.1 for $c=1 / 2$ as follows. Compute the $n$th convergent $p_{n} / q_{n}$ of the continued fraction approximation for $y_{i}$ such that $q_{n} \leq m$ and either $p_{n} / q_{n}=y_{i}$ or $q_{n+1}>m$.
Set $l_{i+1} \leftarrow q_{n} ; y_{i+1} \leftarrow p_{n}$. At this point $\left|l_{i+1} y_{i}-y_{i+1}\right|<1 / \sqrt{d}$.
Compute $x_{i+1} \in \mathbf{Z}$ such that

$$
\left|\frac{x_{i}-\left(d y_{i+1} \bmod 2\right)}{2}-x_{i+1}\right| \leq \frac{1}{2}
$$

Set $\gamma_{i+1} \leftarrow x_{i+1}+\left(d y_{i+1} \bmod 2+y_{i+1} \sqrt{d}\right) / 2$.
Set $\rho_{i+1} \leftarrow l_{i+1} \rho_{i-1}-\gamma_{i+1} \rho_{i} ; \Pi \leftarrow \Pi \cup\left\{\pi_{j} \mid 1 \leq j \leq k, e_{l, j} \geq 1\right\}$.
At this point $\left|\mathbf{N} \rho_{i+1}\right|<\left|\mathbf{N} \rho_{i}\right| / 2$.
IF $d>0$ and $\left|\mathbf{N} \rho_{i+1}\right| \neq 1$ THEN:
We adjust $\rho_{i+1}$ such that $R \rho_{i+1}$ and $I \rho_{i+1}$ do not become too large as follows:
Compute $k \geq 0$ such that $\left|\phi\left(\varepsilon_{1}^{k}\right)-\phi\left(\rho_{i+1}\right)\right|<\phi\left(\varepsilon_{1}\right) / 2$.
Set $\rho_{i+1} \leftarrow \rho_{i+1} / \varepsilon_{1}^{k}$. At this point $\left|R \rho_{i+1}\right|$ and $\left|I \rho_{i+1}\right|$ are both $<\sqrt{\left|\mathbf{N} \rho_{i+1}\right|}$. $\sqrt{\varepsilon_{1}}$. (Refer to Section 2 for an explanation of these facts.)
Step 3: Remove extraneous factors from $\rho_{i-1}$ introduced by the $l_{i}$.
Set $\beta \leftarrow 1$. FOR $\pi \in \Pi$ DO
Compute the maximal $e, f$ such that $\pi^{e}\left|\xi, \pi^{f}\right| \eta$.
Set $\beta \leftarrow \beta \times \pi^{\min (e, f)} ; \rho_{i-1} \leftarrow \rho_{i-1} / \pi^{\min (e, f)}$.
WHILE $\pi \mid \rho_{i-1}$ DO $\rho_{i-1} \leftarrow \rho_{i-1} / \pi$.
RETURN $\delta \leftarrow \beta \rho_{i-1}$.
If one also needs the extended Euclidean scheme $\sigma \xi+\tau \eta=\delta$, it suffices to set $\sigma=$ $\left(s_{1}+\sqrt{d} s_{2}\right) / 2, \tau=\left(t_{1}+\sqrt{d} t_{2}\right) / 2, s_{1}-d s_{2}=2 s_{3}, t_{1}-d t_{2}=2 t_{3}$, and solve the four resulting linear equations with integer coefficients in the integers $s_{1}, s_{2}, s_{3}, t_{1}, t_{2}, t_{3}$ by some integer linear system solver, see Kannan and Bachem [13].

It is easy to show that for fixed $d$ this algorithm has time complexity $O\left(S^{3}\right)$, where $S$ is the size of the input, defined in Section 1. Some additional effort is needed to show that the complexity can be improved to $O\left(S^{2}\right)$, and lastly to determine the complexity in dependence of $d$.

THEOREM 4.1. After input-independent preconditioning, the $G C D$ of $\xi, \eta \in$ $O_{d}, O_{d}$ a UFD, can be computed in

$$
O(S \sqrt{|d|}(\log d)(S+\sqrt{|d|} \log d)), \quad S=\operatorname{size} \xi+\operatorname{size} \eta
$$

binary steps, using classical integer arithmetic procedures.
Proof. This can be shown similarly as for the Euclidean algorithm applied to rational integers (see Knuth [14, Exercise 30 of $\S 4.5 .2]$ ), but the details are somewhat more involved. We have to analyze the preconditioned algorithm, and also to apply
some minor modifications. In what follows, we mean by 'constant' some quantity that is independent of $\xi$ and $\eta$, though it may still depend on $d$. First we have to consider the number of iterations of Step 2. It follows from $\left|\mathbf{N} \rho_{i+1}\right|<\left|\mathbf{N} \rho_{i}\right| / 2$ that this number is bounded by $\log \left|\mathbf{N} \rho_{1}\right|+1$, which is $O(S)$.

Second, we have to establish an upper bound for the size of the remainders $\rho_{2}, \rho_{3}$, etc. Because of the adjustment of $\rho_{i+1}$ at the end of each iteration of Step 2, applied for $d>0$, there exists a constant $C_{1}$ such that both $\left|R\left(\rho_{i+1}\right)\right|$ and $\left|I\left(\rho_{i+1}\right)\right| \leq C_{1} \sqrt{\left|\mathbf{N} \rho_{i+1}\right|}<C_{1} \sqrt{\left|\mathbf{N} \rho_{1}\right|}$, where $C_{1}=\sqrt{\varepsilon_{1}}$ (see Section 2). Thus, size $\rho_{i}=O\left(S+\log \varepsilon_{1}\right)$, without fixing $i$. Moreover, $\varepsilon_{1}<d^{\sqrt{|d|}}$ (cf. Hua [10]), thus $\log \left(\varepsilon_{1}\right)=O(\sqrt{|d|} \log |d|) .{ }^{* * *}$ We will also need the constant $C_{1}$ in the following paragraphs.

It is now easy to analyze the complexity of Step 3 . We assume that the multiplications $\beta \leftarrow \beta \cdot \pi^{\min (e, f)}$ and divisions $\rho_{i-1} \leftarrow \rho_{i-1} / \pi^{\min (e, f)}$ are replaced by sequences of multiplications $\beta \leftarrow \beta \cdot \pi$ and divisions $\rho_{i-1} \leftarrow \rho_{i-1} / \pi$, which does not reduce the complexity. Then, apart from the final multiplication $\delta \leftarrow \beta \rho_{i-1}$, Step 3 consists of a sequence of divisions and multiplications, where one of the operands has size $O(S+\sqrt{|d|} \log d)$ while the other is a prime factor of a rational integer $l<\sqrt{|d|}$. Such a prime factor is either a rational prime $p<\sqrt{|d|}$, or an element $\pi \in O_{d}$ with $|N \pi|<\sqrt{|d|}$. In the latter case, $\pi$ will also have been adjusted by multiplication by an appropriate power of $\varepsilon_{1}$, so that $\operatorname{size}(\pi)=O\left(\log |d|+\log \varepsilon_{1}\right)=O(\sqrt{|d|} \log d)$. Thus each operation requires $O((S+\sqrt{|d|} \log d) \sqrt{|d|} \log d)$ steps. Notice that whenever a division $\alpha / \psi$ is performed, $\mathbf{N} \psi$ must be computed, which takes only $O\left(d(\log d)^{2}\right)$ binary steps, since $\psi$ will always be the second type of operator, with size $O(\sqrt{|d|} \log d)$. The number of multiplications and divisions to be performed in Step 3 is easily seen to be $O(z)$, where $z$ is the sum of the number of prime factors of $\xi, \eta$, and $\rho_{i-1}$, which is $O(S)$. In the final multiplication $\delta \leftarrow \beta \rho_{i-1}$ both factors have once again size $O(S+\sqrt{|d|} \log d)$. Thus the total complexity of all operations in Step 3 becomes
$O\left(S(S+\sqrt{|d|} \log d) \sqrt{|d|} \log d+(S+\sqrt{|d|} \log d)^{2}\right)=O(S(S+\sqrt{|d|} \log d) \sqrt{|d|} \log d)$.
It remains to analyze the complexity of Step 2.
If we compute the exact representation of $x_{i}$ and $y_{i}$, then this calculation would take $O\left(S^{2}\right)$ steps for a fixed $d$, which is too much to guarantee the desired overall complexity for all iterations. However, it suffices to compute floating-point approximations for $x_{i}, y_{i}$, with an appropriate upper bound for the absolute error, 0.01 say. Then instead of $\left|\mathbf{N} \rho_{i+1}\right|<\left|\mathbf{N} \rho_{i}\right| / 2$ we will accomplish $\left|\mathbf{N} \rho_{i+1}\right|<C_{2}\left|\mathbf{N} \rho_{i}\right|$ for some constant $C_{2}, \frac{1}{2}<C_{2}<1$. This does not affect the analysis as far as asymptotic step count and bit sizes are concerned. Notice that the floating-point approximations are thus local to Step 2, and no error gets propagated to the next iteration. The exact values of $x_{i}, y_{i}$ are

$$
x_{i}=2 \frac{R \rho_{i-1} R \rho_{i} \pm I \rho_{i-1} I \rho_{i}}{\mathbf{N} \rho_{i}}
$$

[^1]with the + sign applying for $d<0$, and
$$
y_{i}=2 \frac{-R \rho_{i-1} I \rho_{i}+I \rho_{i-1} R \rho_{i}}{\mathbf{N} \rho_{i} \sqrt{|d|}} .
$$

Let $A_{i}=\log \sqrt{\left|\mathbf{N} \rho_{i-1}\right|}-\log \sqrt{\left|\mathbf{N} \rho_{i}\right|}$. The numerators of $x_{i}$ and $y_{i}$ have absolute values $\leq 2 C_{1}^{2} \sqrt{\left|\mathbf{N} \rho_{i-1}\right|\left|\mathbf{N} \rho_{i}\right|}$. Hence, there is a constant $C_{3}$ such that $\log \left|x_{i}\right|$, $\log \left|y_{i}\right| \leq A_{i}+C_{3}$, and a constant $C_{4}$ such that the desired accuracy for $x_{i}, y_{i}$ is guaranteed, if $A_{i}+C_{4}$ significant digits of both numerator and denominator of $x_{i}$ and $y_{i}$ are computed, and floating-point division is applied.

For the following sufficiently accurate error analysis we need to introduce a variant of the definition of $\operatorname{size}(\alpha)\left(\alpha \in O_{d}\right)$ given in the introduction. (The previous definition is insufficient because size $(a+a \sqrt{d})$ can be roughly $2 \operatorname{size}(a)(a \in \mathbf{Z})$, whereas we will need a bound for $\operatorname{size}(a+a \sqrt{d})-\operatorname{size}(a)$ that is independent of $a$.) Define

$$
\operatorname{size}_{1}(\alpha)=\max (\log |R \alpha|, \log |I \alpha|)
$$

Now the relation $\left|R\left(\rho_{j}\right)\right|,\left|I\left(\rho_{j}\right)\right| \leq C_{1} \sqrt{\left|\mathbf{N} \rho_{j}\right|}$ also allows us to find another constant $C_{5}$ such that

$$
\begin{equation*}
\left|\operatorname{size}_{1}\left(\rho_{j}\right)-\log \sqrt{\left|\mathbf{N} \rho_{j}\right|}\right| \leq C_{5} . \tag{4.2}
\end{equation*}
$$

Therefore, at most $2 C_{5}$ leading digits can cancel out when the subtraction in the calculation of $\mathbf{N} \rho_{j}$ is performed. Applying (4.2) to $j=i-1$ and to $j=i, A_{i}$ can be estimated by ( $\operatorname{size}_{1} \rho_{i-1}-\operatorname{size}_{1} \rho_{i}$ ) within an error bounded by a constant bound $C_{6}$; so the number of digits that ultimately have to be used is ( $\operatorname{size}_{1} \rho_{i-1}-$ size $\left._{1} \rho_{i}\right)+C_{4}+2 C_{5}+C_{6}=B$. Note that $B$, contrary to $A_{i}$, is known before $x_{i}, y_{i}$ are computed. As pointed out, $C_{1}=O\left(\sqrt{\varepsilon_{1}}\right)$. Then it follows that $C_{3}, \ldots, C_{6}$ are all of size $O\left(\log \varepsilon_{1}\right)$, which is $O(\sqrt{|d|} \log d)$, as stated above. Hence $B=$ $O\left(A_{i}+\sqrt{|d|} \log d\right)$. All arithmetic with $B$-digit precision can be performed in $O\left(B^{2}\right)=O\left(\left(A_{i}+\sqrt{|d|} \log d\right)^{2}\right)=O\left(A_{i}^{2}+|d|(\log d)^{2}\right)$ steps. The total complexity of all these operations for all iterations of Step 2 is therefore

$$
O\left(\sum_{i \geq 1} A_{i}^{2}+S|d|(\log d)^{2}\right)=O\left(S\left(S+|d|(\log d)^{2}\right)\right)
$$

Floating-point arithmetic only applies to $x_{i}$ and $y_{i}$, not to other intermediate results like $p_{j}, q_{j}(j=1, \ldots, n), x_{i+1}$ etc. We now consider the complexity of the remaining operations in an iteration of Step 2. The continued fraction approximation $p_{n} / q_{n}$ of $y_{i}$ can be calculated as follows: from the $B$-digit floating-point representation of $y_{i}$ we obtain a fractional representation $s / t, s, t \in \mathbf{Z}$, where $s$ has $B$ digits and $t$ is a power of 2 . Note that $\left|y_{i}\right|$ cannot be larger than its $B$-digit mantissa $s$, since $B$ was large enough to guarantee an absolute error $<1$. If $\left|y_{i}\right|<1 / \sqrt{|d|}$, then instead of computing continued fractions we may simply choose $l_{i+1}=1, y_{i+1}=0$, which satisfies the required inequality $\left|l_{i+1} y_{i}-y_{i+1}\right|<1 / \sqrt{|d|}$. If not, then both $\log (s)$ and $\log (t)$ are $O\left(A_{i}+\sqrt{|d|} \log d\right)$. If $g_{1}, g_{2}$ are the quotients of the Euclidean algorithm for rational integers with nonnegative remainders, applied to $s, t$, then
the numbers $p_{i}, q_{i}$ are determined by the following recurrence relation:

$$
\left.\begin{array}{c}
p_{-1}=q_{0}=0, \quad q_{-1}=p_{0}=1, \\
p_{i}=p_{i-2}+g_{i} p_{i-1} \\
q_{i}=q_{i-2}+g_{i} q_{i-1}
\end{array}\right\}, \quad i=1,2, \ldots .
$$

From these equations it becomes clear that

$$
\sum_{i=1}^{n} \log \left(g_{i}\right)=O\left(\log q_{n}\right)=O(\log \sqrt{|d|})
$$

and since the intermediate results of the Euclidean algorithm applied to $s, t$, as well as the numbers $p_{i}, q_{i}$ have $O\left(A_{i}+\sqrt{|d|} \log d\right)$ digits, we find by the usual method that $O\left((\log d)\left(A_{i}+\sqrt{|d|} \log d\right)\right)$ binary steps are needed for the calculation of $p_{n}, q_{n}$.

The logarithms of the rational and irrational part of $\gamma_{i+1}$ are $O\left(A_{i}+\sqrt{|d|} \log d\right)$, since $\gamma_{i+1}$ is an approximation of $l_{i+1} \rho_{i-1} / \rho_{i}$. This ensures the desired complexity bound $O\left(\left(A_{i}+\sqrt{|d|} \log d\right)(S+\sqrt{|d|} \log d)\right)$ for the multiplication $\gamma_{i+1} \rho_{i}$ in the computation of $\rho_{i+1}$. Then the same bound for the order of magnitude applies to the logarithms of $x_{i+1}, y_{i+1}, p_{1}, \ldots, p_{n}$. The only step which requires further consideration is the division of $\rho_{i+1}$ by an appropriate power of $\varepsilon_{1}$.

Consider first the value $\rho_{i+1}^{(0)}$ of $\rho_{i+1}$ before this adjustment. Recall that $\rho_{i+1}^{(0)}=$ $\rho_{i} \tau$, where $|R \tau|,|I \tau|<1$ (actually $\leq \frac{1}{2}$ if the exact values of $x_{i}, y_{i}$ are used). This implies $\operatorname{size}_{1} \rho_{i+1}^{(0)} \leq \operatorname{size}_{1} \rho_{i}+C_{7}$ for some constant $C_{7}$ with $C_{7}=O(\log |d|)$. Let $\rho_{i+1}^{(t)}$ be the value of $\rho_{i+1}$ after the adjustment. Then $\phi\left(\rho_{i+1}^{(0)}\right)=\phi\left(\rho_{i+1}^{(t)}\right)+k \phi\left(\varepsilon_{1}\right) \geq$ $(k-1) \phi\left(\varepsilon_{1}\right)$. Together with the defining equality for the argument function $\phi$, this implies $k=O\left(A_{i}\right)$. We can proceed as follows: first, if $R \rho_{i+1}^{(0)}<0$, replace $\rho_{i+1}^{(0)}$ by $-\rho_{i+1}^{(0)}$. Then, if $I \rho_{i+1}^{(0)}>0$, compute $\rho_{i+1}^{(j)}=\rho_{i+1}^{(j-1)} / \varepsilon_{1}, j=1,2, \ldots$ until $I \rho_{i+1}^{\left(j_{0}\right)} \leq 0$, so that $\phi \rho_{i+1}^{\left(j_{0}-1\right)}>0, \phi \rho_{i+1}^{\left(j_{0}\right)} \leq 0$, then choose among $\rho_{i+1}^{\left(j_{0}-1\right)}$ and $\rho_{i+1}^{\left(j_{0}\right)}$ the one for which the rational part, or, equivalently, the absolute value of the irrational part, is smaller. If instead after the possible replacement of $\rho_{i+1}^{(0)}$ by $-\rho_{i+1}^{(0)}$ we get $I \rho_{i+1}^{(0)}<0$, compute similarly $\rho_{i+1}^{(j)}=\rho_{i+1}^{(j-1)} \varepsilon_{1}, j=1,2, \ldots$, until $I \rho_{i+1}^{\left(j_{0}\right)} \geq 0$. The number of divisions is $O\left(A_{i}\right)$, and each division has complexity $O((S+\sqrt{|d|} \log d) \sqrt{|d|} \log d)$.

We conclude that the highest asymptotic complexity within Step 2 occurs for the adjustment of $\rho_{i+1}$, namely $O\left(A_{i}(S+\sqrt{|d|} \log d) \sqrt{|d|} \log d\right)$. Adding this up for $i=1,2, \ldots$, we obtain $O(S(S+\sqrt{|d|} \log d) \sqrt{|d|} \log d)$. We have established that this is indeed an upper bound for the complexity of all operations, including Step 3. This completes the proof.

Omitting the adjustment of $\rho_{i+1}$ does not affect the correctness of the algorithm. In this case, it is not obvious, however, that the complexity $O\left(S^{2}\right)$ can still be achieved for fixed $d$. Moreover, it is in any case desirable to carry out this normalization at least for the final result. After all, if $\operatorname{GCD}(\xi, \eta)=1$, it would be unsatisfactory to obtain some power $\pm \varepsilon_{1}^{k}, k \neq 0$, as output.

We now discuss how the preconditioning for this algorithm can be done. For the calculation of the fundamental unit $\varepsilon_{1}$ in the case $d>0$, there exists the well-known algorithm using continued fractions. An improvement has been provided in Pohst and Zassenhaus [20]. The factorization of the multipliers $l<\sqrt{|d|}$ will begin by finding their factorization into rational primes $p$. Using the sieve of Eratosthenes,
the largest prime factor for every $l$ can be found in time $O(\sqrt{|d|})$, which essentially solves the problem. Then it can be determined efficiently, using the law of quadratic reciprocity, which of these rational primes $p$ split further in $O_{d}$; see Section 5 . One of the referees suggests the following algorithm for finding the prime factors $\pi_{1}$ and $\pi_{2}$, if $p$ splits further. The idea is to explicitly find the transformation from the quadratic form $(1,0,-d / 4)$ or $(1,1,(1-d) / 4)$, whichever is integral, to a form $\left(4 p, g_{1}, g_{2}\right)$, all of discriminant $d$. Consider first the case $d \equiv 0(\bmod 4)$. Then the transforming matrix $M=\left[\begin{array}{cc}\alpha \\ \beta & \gamma \\ \delta\end{array}\right]$ yields the solution $\alpha^{2}-d \beta^{2}=4 p$, the wanted factorization. For $d \equiv 1(\bmod 4)$, the matrix $M$ yields $(2 \alpha+\beta)^{2}-d \beta^{2}=4 p$ as the solution. In order to compute $M$, one appeals to the theory of reduction on the principal cycle of reduced quadratic forms of discriminant $d$ (cf. [25, Section 4, pp. 256-261]). One starts at a form $\left(4 p, g_{1}, g_{2}\right)$ or $\left(p, g_{1}, g_{2}\right)$ of discriminant $d$, which is easily determined, using an algorithm for taking square roots modulo $p$. Then one steps through the principal cycle by repeated reduction until one detects the forms $(1,0,-d / 4)$ or $(1,1,(1-d) / 4)$, respectively. For each prime $p$, the running time depends on the number of forms in the cycle, about $O\left(d^{1 / 2+\varepsilon}\right)$, such that the procedure for all primes $<\sqrt{|d|}$ takes $O\left(d^{1+\varepsilon}\right)$.

We conclude by briefly discussing the application of this algorithm to domains $O_{d}$ with $h>1$. The algorithm itself does not change, except that prime ideal factors of ( $l$ ) for the multipliers $l$ have to be considered, rather than prime factors of $l$ in $O_{d}$. What does it accomplish? Of course, it would be a trivial task to find a representation for the greatest common divisor of the ideal $(\xi)$ and $(\eta)$, as $(\xi, \eta)$ would already be such a representation. Instead, it is desirable to put it into some normal form. Let $\left\{I_{1}, \ldots, I_{h}\right\}$ be a set of representatives of the ideal classes. As is well known, the $I_{i}$ can be chosen such that

$$
\left\|I_{i}\right\| \leq \begin{cases}\frac{\sqrt{d}}{2} & \text { for } d>0 \\ \frac{2 \sqrt{|d|}}{\pi} & \text { for } d<0\end{cases}
$$

where $\|I\|$ is the number of congruence classes $\bmod I$. Then every ideal of $O_{d}$ has a unique representation of the form $(\alpha) / I_{i}, \alpha \in O_{d}, 1 \leq i \leq h$. The GCD algorithm can now be used to compute the normal form of $(\xi, \eta)$. The following preparatory work is needed: one has to compute the normal forms of the ideal prime factors of $l, l=2, \ldots,\lfloor\sqrt{|d|}\rfloor$, and a multiplication table for the normal forms of the ideal products $I_{i} \cdot I_{j}$. The latter problem is equivalent to the computation of the class group and can be done in time $|d|^{1 / 4+o(1)}$, see Schoof [25]. The factorization of the ideals $(l)$ is at least as difficult as in the case $h(d)=1$. After this preparatory work has been finished, the complexity bound of Theorem 4.1 applies again.
5. A Constructive Version of a Theorem by A. Thue with Application. A theorem by Axel Thue [30] states that if $a$ and $m$ are positive relatively prime integers, then there exist integers $0<x \leq \sqrt{m}, 0<|y| \leq \sqrt{m}$ such that $a x+y \equiv 0 \bmod m$. This theorem and its generalizations (cf. Brauer and Reynolds [2] and Nagell [19]) are usually proved using the pigeon hole principle. The following theorem shows how all solutions for the above congruence can be found in $O\left(\log ^{2} m\right)$ steps.

THEOREM 5.1. Let $a \geq 1, m, e, f \geq 2$ be integers such that $a<m$, $(e-1)$. $(f-1)<m<e f$. Then the problem

$$
\begin{equation*}
m|a x+y, \quad 0<x<e,|y|<f, y \neq 0 \tag{5.1}
\end{equation*}
$$

is solvable in integers $x, y$ if and only if $d=\operatorname{GCD}(a, m)<f$. Furthermore, assuming that this is the case, let

$$
\frac{p_{0}}{q_{0}}=\frac{0}{1}, \frac{p_{1}}{q_{1}}, \ldots, \frac{p_{N}}{q_{N}}=\frac{a / d}{m / d}, \quad q_{N} \geq \frac{m}{f-1}>e-1,
$$

be the continued fraction approximations of a/m, and choose $n$ such that $q_{n}<e \leq$ $q_{n+1}$. Then $x_{1}=q_{n}, y_{1}=m p_{n}-a q_{n}$ is a solution for (5.1). The set of all solutions for (5.1) exclusively either consists of $\lambda x_{1}, \lambda y_{1}, 1 \leq \lambda<\min \left(e / x_{1}, f /\left|y_{1}\right|\right)$, or else consists of $x_{1}, y_{1}$ and $x_{2}, y_{2}$ with $y_{1} y_{2}<0$. In the latter case, we can determine $x_{2}, y_{2}$ from $p_{n-1} / q_{n-1}$ or $p_{n+1} / q_{n+1}$ in $O\left(\log ^{2} m\right)$ steps.

Before we can prove this theorem, we need to establish a lemma from the theory of continued fraction approximations. Following Hardy and Wright [8, Section 10], we denote a continued fraction by

$$
\left[a_{0}, a_{1}, \ldots, a_{n}, \ldots\right]=a_{0}+\frac{1}{\left[a_{1}, \ldots, a_{n}, \ldots\right]}
$$

The $n$th convergent is given by

$$
\frac{p_{n}}{q_{n}}=\left[a_{0}, a_{1}, \ldots, a_{n}\right]
$$

and satisfies $p_{n} q_{n-1}-p_{n-1} q_{n}=(-1)^{n-1}$. Notice that

$$
\begin{equation*}
\left[a_{0}, a_{1}, \ldots, a_{n}, 1\right]=\left[a_{0}, a_{1}, \ldots, a_{n}+1\right] \tag{5.2}
\end{equation*}
$$

but this is the only amibiguity possible for the simple continued fraction expansion of a real number where $a_{i}, i>0$, are positive integers.

Lemma 5.1 ([8, Theorem 172]). If

$$
x=\frac{P Z+R}{Q Z+S}
$$

where $Z \geq 1$ and $P, Q, R$, and $S$ are integers such that

$$
Q>S>0, \quad P S-Q R= \pm 1
$$

then $R / S$ and $P / Q$ are two consecutive convergents in the simple continued fraction expansion of $x$, provided we choose the left-hand side of (5.2) to resolve ambiguity.

Proof. Consider the continued fraction expansion

$$
\frac{P}{Q}=\left[a_{0}, a_{1}, \ldots, a_{n}\right]=\frac{p_{n}}{q_{n}}
$$

From (5.2) it follows that we may choose $n$ even or odd as we need it. Now let $n$ be such that

$$
P S-Q R=(-1)^{n-1}=p_{n} q_{n-1}-p_{n-1} q_{n}
$$

Now $\operatorname{GCD}(P, Q)=1$ and $Q>0$, hence $P=p_{n}$ and $Q=q_{n}$, and therefore

$$
p_{n}\left(S-q_{n-1}\right)=q_{n}\left(R-p_{n-1}\right) .
$$

This implies that $q_{n} \mid\left(S-q_{n-1}\right)$, which by virtue of $q_{n}=Q>S>0$ and $q_{n}>q_{n-1}>0$ is only possible if $S-q_{n-1}=0$. Therefore, $R=p_{n-1}$ and $S=q_{n-1}$, and

$$
x=\frac{p_{n} Z+p_{n-1}}{q_{n} Z+q_{n-1}} \quad \text { implies } \quad x=\left[a_{0}, a_{1}, \ldots, a_{n}, Z\right]
$$

Proof of Theorem 5.1. Since $m \mid a x+y$, there exists an integer $z$ such that $y=m z-a x$ and thus $d \mid y$, which implies $d \leq|y|<f$. For $p_{n} / q_{n}$ we have, as in the proof of Lemma 3.1,

$$
\left|\frac{a}{m}-\frac{p_{n}}{q_{n}}\right| \leq \frac{1}{q_{n} q_{n+1}} \leq \frac{1}{q_{n} e}
$$

therefore, $\left|y_{1}\right|=\left|a q_{n}-m p_{n}\right| \leq m / e<f$. Notice that $y_{1} \neq 0$ because $a / m \neq p_{n} / q_{n}$ for $q_{n}<q_{N}=m / d$.

We prove next that if $\hat{x}_{1}, \hat{y}_{1}$ also solve (5.1), and $y_{1} \hat{y}_{1}>0$, then $y_{1} / x_{1}=\hat{y}_{1} / \hat{x}_{1}$. Since $\hat{x}_{1}\left(a x_{1}+y_{1}\right) \equiv x_{1}\left(a \hat{x}_{1}+\hat{y}_{1}\right) \equiv 0 \bmod m$, we have $y_{1} \hat{x}_{1}-\hat{y}_{1} x_{1} \equiv 0 \bmod m$. But $\left|y_{1} \hat{x}_{1}\right|,\left|\hat{y}_{1} x_{1}\right| \leq(e-1)(f-1)<m$, hence $y_{1} \hat{x}_{1}=\hat{y}_{1} x_{1}$, which is our claim. Nonetheless, it can happen that $\operatorname{GCD}\left(x_{1}, y_{1}\right)=g \neq 1$, and therefore the assertion $\hat{x}_{1}=\lambda x_{1}, \hat{y}_{1}=\lambda y_{1}$ needs proof. First we note that

$$
a x_{1}+y_{1}=m p_{n} \quad \text { and } \quad \operatorname{GCD}\left(p_{n}, x_{1}\right)=\operatorname{GCD}\left(p_{n}, q_{n}\right)=1
$$

Thus,

$$
a \lambda \frac{x_{1}}{g}+\lambda \frac{y_{1}}{g}=\frac{\lambda m p_{n}}{g} \equiv 0 \bmod m
$$

if and only if $g \mid \lambda$, since $\operatorname{GCD}\left(g, p_{n}\right)=1$.
Assume now that there exists a second solution $x_{2}, y_{2}$ to (5.1) such that $y_{1} y_{2}<0$. Let $a x_{1}+y_{1}=m z_{1}, a x_{2}+y_{2}=m z_{2}$. Again, by multiplying the first equation with $x_{2}$, the second with $x_{1}$ and subtracting, we get $x_{1} y_{2}-x_{2} y_{1} \equiv 0 \bmod m$. Since $\left|x_{1} y_{2}\right|,\left|x_{2} y_{1}\right|<m$ and $y_{1} y_{2}<0$, we must have $\left|x_{1} y_{2}-x_{2} y_{1}\right|=m$. Thus, $m\left|x_{1} z_{2}-x_{2} z_{1}\right|=\left|x_{1} y_{2}-x_{2} y_{1}\right|=m$, which implies that

$$
\begin{equation*}
\left|x_{1} z_{2}-x_{2} z_{1}\right|=1 \tag{5.3}
\end{equation*}
$$

One immediate conclusion from (5.3) is that no solutions proportional to either $x_{1}, y_{1}$ or $x_{2}, y_{2}$, the only other possible solutions (as shown before), can occur. For otherwise, for any of these solutions, say $\hat{x}, \hat{y}$, one has $\operatorname{GCD}(\hat{x}, \hat{z}) \neq 1$, where $\hat{z}$ is the corresponding, also proportional, multiplier of $m$.

It is harder to show how this second solution $x_{2}, y_{2}$ can be computed in case it exists. Surprisingly, this alternate solution can arise in two different ways. Without loss of generality, let us assume that $x_{2}<x_{1}$. Notice that by this assumption we now only know that either $x_{1}, y_{1}$ or $x_{2}, y_{2}$ is the solution found as stated in the theorem.

Case 1: $\left|y_{2}\right|>\left|y_{1}\right|$. Let $Z=\left|y_{2} / y_{1}\right|>1$. Then $Z m z_{1}=Z a x_{1}+Z y_{1}=Z a x_{1}-y_{2}$, hence $Z m z_{1}+m z_{2}=Z a x_{1}+a x_{2}$, or

$$
\frac{a}{m}=\frac{Z z_{1}+z_{2}}{Z x_{1}+x_{2}} .
$$

All conditions to Lemma 5.1 with $P=z_{1}, Q=x_{1}, R=z_{2}$, and $S=x_{2}$ are now satisfied (refer in particular to (5.3)), and we can conclude that $z_{2} / x_{2}$ and
$z_{1} / x_{1}$ must be consecutive convergents to $a / m$. Therefore, $x_{2}=q_{n-1}$ and $y_{2}=$ $m p_{n-1}-a q_{n-1}$.

Case 2: $\left|y_{2}\right| \leq\left|y_{1}\right|$. Consider for an integer $k \geq 0$,

$$
x_{3}=x_{1}+k x_{2}, \quad y_{3}=y_{1}+k y_{2}, \quad z_{3}=z_{1}+k z_{2},
$$

such that

$$
\left|y_{2}\right| \geq\left|y_{3}\right| \quad \text { and } \quad \operatorname{sign}\left(y_{3}\right)=-\operatorname{sign}\left(y_{2}\right)
$$

Now

$$
0<x_{2}<x_{3}, \quad \frac{a}{m}=\frac{Z z_{3}+z_{2}}{Z x_{3}+x_{2}}, \quad Z=\left|\frac{y_{2}}{y_{3}}\right| \geq 1, \quad z_{3} x_{2}-z_{2} x_{3}= \pm 1
$$

and Lemma 5.1 applies again. Therefore, $z_{2} / x_{2}$ and $z_{3} / x_{3}$ are consecutive continued fraction approximations. Notice that $Z=1$ is possible if and only if $z_{2} / x_{2}$ is the second to the last convergent of $a / m$. In that case,

$$
\frac{z_{3}}{x_{3}}=\left[0, a_{1}, \ldots, a_{N}-1\right] \quad \text { where } \frac{a}{m}=\left[0, a_{1}, \ldots, a_{N}\right] .
$$

In order to compute $x_{1}$ and $y_{1}$, we must find $k$. Since $\left|y_{1}\right|=\left|y_{3}-k y_{2}\right|$ is monotonically increasing with $k$, we choose the smallest $k$ such that $x_{3}-k x_{2} \leq e-1$. In other words, the only possible value is

$$
k=\left\lceil\frac{x_{3}-e+1}{x_{2}}\right\rceil .
$$

Observe that $x_{1}=x_{3}-k x_{2} \neq x_{2}$, since $\operatorname{GCD}\left(x_{2}, x_{3}\right)=1$, and

$$
a x_{1}+y_{1} \equiv 0 \bmod m \quad \text { for } y_{1}=y_{3}-k x_{2}
$$

The following example shows that all three types of solutions, namely either a single one, or two, or a family of proportional solutions do occur.

Examples. Let $m=11, e=f=4, a=7$. The continued fraction expansion of $7 / 11$ is $[0,1,1,1,3]$ and the convergents are $0 / 1,1 / 1,1 / 2,2 / 3$, and $7 / 11$. Hence, $x_{1}=3, y_{1}=2 \cdot 11-7 \cdot 3=1$, and $x_{2}=2, y_{2}=1 \cdot 11-7 \cdot 2=-3$ are the only solutions for (5.1) in this case. For $a=2$, the only solution is $1 \cdot 2-2 \equiv 0 \bmod 11$, but for $a=1$, there are three solutions, $1 \cdot 1-1 \equiv 2 \cdot 1-2 \equiv 3 \cdot 1-3 \equiv 0 \bmod 11$.

Next, let $m=244, a=47, e=7$, and $f=39$. The first three convergents are $0 / 1,1 / 5$, and $5 / 26$. Thus, $x_{2}=5$ and $y_{2}=244 \cdot 1-47 \cdot 5=9$. The second solution is obtained from $k=\lceil(26-7+1) / 5\rceil=4$ as $x_{1}=26-4 \cdot 5=6$ and $y_{1}=38$.

Finally, let $m=56, a=21, e=7$, and $f=9$. The continued fraction expansion of $21 / 56$ is $[0,2,1,2]$ and the convergents are $0 / 1,1 / 2,1 / 3$, and $3 / 8$. Thus, $x_{2}=3$, $y_{2}=-7$. We obtain $z_{3} / x_{3}=[0,2,1,1]=2 / 5, k=\lceil(5-7+1) / 3\rceil=0, x_{1}=5$, and $y_{1}=7$.

One application of Theorem 5.1, described by Wang [31] and others, is to recover rational numbers from their modular representations. Set $y / x \in \mathbf{Q}$ and suppose we have found bounds $e$ and $f$ for the denominator and numerator, respectively. Then, having computed $a=y x^{-1} \bmod m, \operatorname{GCD}(x, m)=1,(e-1)(f-1)<m$, we can find $y / x$ by continued fraction approximation. Unfortunately, we may get two possible fractions $y_{1} / x_{1}, y_{2} / x_{2}$ of opposite sign. One way to resolve the ambiguity is to choose the $e$ twice the bound of the denominator and select the solution with
$x<e / 2,|y|<f$. Wang, in fact, chooses $m$ such that $\lceil\sqrt{m / 2}\rceil$ is a bound for both numerator and denominator. In this application, the existence of a solution is assumed and Thue's theorem does not come into play.

We now apply Theorem 5.1 to the problem of factoring a rational prime $p$ in the UFD $O_{d}, d$ fixed. We again must precondition our algorithm by factoring all rational primes smaller than $2 \sqrt{|D|}(\sqrt{D}$ suffices for $D>0)$. First consider $p \nmid d$. From Section 2 we know that $p<\sqrt{D}$ factors if and only if $(D / p)=+1$. In that case, a prime factor of $p$ is $\pi=\operatorname{GCD}(p, l+\sqrt{D})$, where $l^{2} \equiv D$ modulo $p$. We can compute $l$ by either the Tonelli-Shanks algorithm (cf. D. Knuth [14, Section 4.6.2, Exercise 15]) or by Schoof's algorithm [26]. The latter is deterministic and runs in $O\left(\log ^{8} p\right)$ steps, since $|d|$ is fixed. The GCD algorithm of Section 3 can give us the wanted factor $\pi$, but in this special case, Theorem 5.1 can be applied to our advantage. Let

$$
e=\lceil\sqrt{p / \sqrt{|D|}}\rceil \quad \text { and } \quad f=\lceil\sqrt{p \sqrt{|D|}}\rceil
$$

Then, $(e-1)(f-1)<p<e f$ and $2 \leq e \leq f$ for $p>\sqrt{|D|}$. Then we compute the continued fraction approximation to $l / p$ and get integers $x, y$ such that $p \mid y l+x$, $0<y<e,|x|<f$. This solution satisfies

$$
0 \equiv(x+l y)(x-l y) \equiv x^{2}-D y^{2} \bmod p
$$

thus there exists an integer $q$ with $x^{2}-D y^{2}=q p$. By our bounds we get

$$
\begin{gathered}
x^{2}-D y^{2}<p \sqrt{|D|}-D \frac{p}{\sqrt{|D|}}=2 \sqrt{|D|} p \text { for } D<0, \\
\left|x^{2}-D y^{2}\right|<\max \left(x^{2}, D y^{2}\right)<\sqrt{D} p \text { for } D>0
\end{gathered}
$$

thus $|q|<2 \sqrt{|D|}$.
In the other case $p \mid d, p>\sqrt{|D|}$, we have $x^{2}-D y^{2}=q p$ with $x=0, y=1$ and $q<\sqrt{D}$. In both cases the factorization of $q$ into primes $q=\gamma_{1} \cdots \gamma_{k}$ in $O_{d}$ is already known. Since $O_{d}$ is a UFD and

$$
(x+\sqrt{D} y)(x-\sqrt{D} y)=\gamma_{1} \cdots \gamma_{k} p
$$

$\gamma_{i}$ must divide $x+\sqrt{D} y$ or $x-\sqrt{D} y$. Let $\gamma$ be a maximum product of $\gamma_{i}$ such that $\gamma$ divides $x+\sqrt{D} y$. Then $q / \gamma$ divides $x-\sqrt{D} y$, and the prime factors of $p$ are then $(x+\sqrt{D} y) / \gamma$ and $(x-\sqrt{D} y) /(q / \gamma)$. To prove this, we only have to show that neither quotient is a unit. This follows from the fact that the division can only decrease the norm of the dividend. If one quotient became a unit, the other one would have to have the norm of their product, $p^{2}$, which is larger than its original norm $q p$.

We now discuss how to factor $\xi \in O_{d}$ with $O_{d}$ a unique factorization domain. We first factor $\mathbf{N} \xi=p_{1} \cdots p_{k}$ over the integers. If $\left(d / p_{i}\right)=+1$ or $p_{i}$ divides $d$, we split $p_{i}=\pi_{i} \bar{\pi}_{i}$ by the álgorithm discussed above. We thus obtain a factorization $\xi \bar{\xi}=\pi_{1} \cdots \pi_{l}$, and it remains to trial-divide $\xi$ by $\pi_{i}, 1 \leq i \leq l$, to determine which are its prime factors.
6. Conclusion. We have described algorithms for taking the greatest common divisor and computing the prime factorization of numbers in quadratic fields with
unique factorization. The methods also apply to computing canonical representations of unions of ideals in quadratic number rings without unique factorization. Our algorithms are of polynomial running time provided we fix the discriminant. We have also shown how to reduce factorization in quadratic number rings with unique factorization to rational integer factorization. If the discriminant is large, say of order $10^{10}$, our algorithms unfortunately become impractical. Future investigations will focus on how to treat these cases efficiently.

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[^1]:    ***Unfortunately, by the Brauer-Siegel theorem [16, Chapter XVI], not a much better estimate for $\log \left(\varepsilon_{1}\right)$ is to be expected, since $h(d) \log \left(\varepsilon_{1}\right) \gg d^{1 / 2-\varepsilon}$ for $d \rightarrow \infty$.

