

Jug measuring: Algorithms and complexity

Min-Zheng Shieh, Shi-Chun Tsai*

Department of Computer Science, National Chiao-Tung University, Hsinchu 30050, Taiwan

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Abstract

We study the hardness of the optimal jug measuring problem. By proving tight lower and upper bounds on the minimum number of measuring steps required, we reduce an inapproximable NP-hard problem (i.e., the shortest GCD multiplier problem [G. Havas, J.-P. Seifert, The Complexity of the Extended GCD Problem, in: LNCS, vol. 1672, Springer, 1999]) to it. It follows that the optimal jug measuring problem is NP-hard and so is the problem of approximating the minimum number of measuring steps within a constant factor. Along the way, we give a polynomial-time approximation algorithm with an exponential error based on the well-known *LLL* basis reduction algorithm.

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1. Introduction

Let α and β be two positive integers. You are given an α -liter jug, a β -liter jug, an unlimited source of water and a drain. You can *fill* a jug full of water from the source, *empty* water in a jug into the drain, or *pour* the water from one jug to another. How do you measure x liters of water? This measuring problem is the so-called *water jug problem*, which has been studied for a long time and is a popular problem for programming contests, a frequent heuristic search exercise in artificial intelligence and algorithms, etc. The water jug problem has several variants, such as the problem of sharing jugs of wine [5], i.e., *given an x -liter jug full of wine and two empty jugs of capacity α and β where $x = \alpha + \beta$, what is the quickest way to divide the wine equally by pouring the wine among the jugs?* This problem can be reduced to the water jug problem, since we can replace “pouring from the x -liter jug to any other jug” by the filling operation, and “pouring from any jug to the x -liter jug” by the emptying operation. In [3], a problem proposed by Ehrlich asks: *for two relatively prime integers α and β and an integer m with $1 \leq m \leq \beta$, is it possible to measure m liters?*

Boldi et al. [4] generalized the water jug problem by considering a set of jugs of fixed capacities and they found out which quantities are *measurable* and proved upper and lower bounds on the number of steps necessary for measuring a specified amount of water. Here a quantity x is called *measurable* if after several steps one of the jugs holds the

* Corresponding author. Tel.: +886 3 5131551.

E-mail addresses: mzhsieh@csie.nctu.edu.tw (M.-Z. Shieh), sctsai@csie.nctu.edu.tw (S.-C. Tsai).

quantity x . In other words, the generalized water jug problem is: *given a set of jugs of fixed capacities, find out which quantities are measurable*. More specifically, suppose that we are given n jugs with positive integral capacities c_i , $i \in [n]$, where $[n]$ denotes the set $\{1, \dots, n\}$.¹ Without loss of generality, we assume that $c_1 \leq c_2 \leq \dots \leq c_n$. Let x be a positive integer $\leq c_n$. Boldi et al. proved that x is measurable if and only if it is a multiple of the greatest common divisor of c_i 's. Define $\mu_c(x) = \min_{\mathbf{x}=\mathbf{x} \cdot \mathbf{c}} \|\mathbf{x}\|_1$, where $\mathbf{x} = (x_1, \dots, x_n) \in \mathbb{Z}^n$ and $\|\mathbf{x}\|_1 = \sum_{i=1}^n |x_i|$. They also proved that, for every measurable x , (1) x can be measured in at most $\frac{5}{2}\mu_c(x)$ steps. (2) No algorithm can measure x in less than $\frac{1}{2}\mu_c(x)$ steps.

In this paper we deal with the *optimal jug measuring problem*, which considers the *minimum* number of measuring steps. We generalize *measurability* by defining that a quantity x is *additively measurable* if there is a sequence of operations such that the sum of the final contents in the jugs is equal to x . We prove that x is additively measurable if and only if it is a multiple of the greatest common divisor of c_i 's. For every additively measurable integer x , we obtain a lower bound $2\mu_c(x) - n_e$, where n_e is the number of non-empty jugs. Consequently, we have that: a measurable integer x cannot be measured in less than $\max\{2\mu_c(x) - n, \mu_c(x)\}$ steps. Furthermore, all the measurable quantities can be measured in at most $2\mu_c(x)$ steps such that the largest jug contains the quantity x and the others are empty. With the tight lower and upper bounds, we can reduce the problem of computing $\mu_c(x)$ to the optimal jug measuring problem. Meanwhile, Havas and Seifert [6] proved that a special case of computing $\mu_c(x)$ is inapproximable in polynomial time within a factor of k , where $k \geq 1$ is an arbitrary constant. This implies that the optimal jug measuring problem is inapproximable in polynomial time within a constant factor. Moreover, we can reduce the problem of computing $\mu_c(x)$ to the closest lattice vector problem. Finally, we propose a polynomial-time approximation algorithm with exponential errors for computing $\mu_c(x)$ based on the famous *LLL* (Lenstra–Lenstra–Lovasz) basis reduction algorithm.

The rest of the paper is organized as follows. In Section 2, we give some notation and definitions. In Section 3, we characterize additive measurability and prove the lower and upper bounds for the number of non-pour operations. In Section 4, we prove the lower and upper bounds for the number of minimum measuring steps. In Section 5, we show that computing $\mu_c(x)$ is inapproximable and give a polynomial-time approximation algorithm with exponential errors. Inapproximable results on the jug measuring problem are also addressed. Section 6 concludes the paper.

2. Notation and definitions

Following Boldi et al. [4], we define three types of elementary operations on the jugs.

Definition 1. Elementary jug operations:

- (1) $\downarrow i$: *fill* the i th jug (from the source) up to its capacity, and we call it the *fill* operation;
- (2) $i \uparrow$: *empty* the i th jug (into the drain) completely, and we call it the *empty* operation;
- (3) $i \rightarrow j$: *pour* the contents of the i th jug to the j th jug, $i \neq j$, and we call it the *pour* operation. Note that pour operation never changes the total sum of the contents, and at the end of this operation, the i th jug is empty or the j th jug is full.

Let O denote the set of all possible elementary operations, that is, $O = \{\downarrow i \mid \forall i \in [n]\} \cup \{i \uparrow \mid \forall i \in [n]\} \cup \{i \rightarrow j \mid \forall i, j \in [n], i \neq j\}$. We say $\sigma \in O^*$ is a sequence of operations, and use $|\sigma|$ to denote the length of σ (number of operations in σ). We use ϵ to denote the empty sequence, i.e., $|\epsilon| = 0$. Let \mathbb{N} be the set of non-negative integers. A *state* is a vector $\mathbf{s} \in \mathbb{N}^n$, where s_i denotes the amount contained in jug i . We define the state-transition function as follows.

Definition 2. A function $\delta : \mathbb{N}^n \times O^* \rightarrow \mathbb{N}^n$ is a state-transition function if:

- (1) $\delta(\mathbf{s}, \epsilon) = \mathbf{s}$;
- (2) $\delta(\mathbf{s}, \downarrow i) = (s_1, \dots, s_{i-1}, c_i, s_{i+1}, \dots, s_n)$;
- (3) $\delta(\mathbf{s}, i \uparrow) = (s_1, \dots, s_{i-1}, 0, s_{i+1}, \dots, s_n)$;
- (4) $\delta(\mathbf{s}, i \rightarrow j) = (t_1, \dots, t_n)$, where $t_k = s_k$ for all $k \notin \{i, j\}$, $t_i = \max\{0, s_i - (c_j - s_j)\}$, and $t_j = \min\{c_j, s_i + s_j\}$;
- (5) for $|\sigma| \geq 1$ and $o \in O$, $\delta(\mathbf{s}, \sigma o) = \delta(\delta(\mathbf{s}, \sigma), o)$.

¹ For convenience we use $[n]$ to denote $\{1, 2, \dots, n\}$, although $[n]$ stands for $\{0, 1, \dots, n-1\}$ under von Neumann's definition of natural numbers.

We say a state s is *reachable* if $\delta(\mathbf{0}, \sigma) = s$ by some sequence $\sigma \in O^*$. Furthermore, we say σ is *optimal* if it is the shortest one that reaches s . $\delta_i(s, \sigma)$ denotes the i -th entry of $\delta(s, \sigma)$. $e_i(\sigma)$ denotes the number of $i \uparrow$ operations in σ and $e(\sigma) = \sum_{i \in [n]} e_i(\sigma)$. $f_i(\sigma)$ denotes the number of $\downarrow i$ operations in σ and $f(\sigma) = \sum_{i \in [n]} f_i(\sigma)$. $p_i(\sigma)$ denotes the number of pour operations applied to jug i and $p(\sigma) = \frac{1}{2} \sum_{i \in [n]} p_i(\sigma)$. A quantity $x \in \mathbb{N}$ is measurable via sequence σ iff there exists $i \in [n]$ such that $\delta_i(\mathbf{0}, \sigma) = x$. For convenience, let $\mathbf{c} = \{c_1, \dots, c_n\}$ and $\gcd(\mathbf{c})$ denote the greatest common divisor of c_1, \dots, c_n . The set of quantities that are measurable using the capacities in \mathbf{c} is denoted as $\mathbf{M}(\mathbf{c})$. Boldi et al. [4] proved that all of the measurable quantities are multiples of the greatest common divisor of the capacities as stated in the following.

Proposition 1 ([4]). $\mathbf{M}(\mathbf{c}) = \{m \cdot \gcd(\mathbf{c}) \mid \text{for all non-negative integer } m \leq c_n / \gcd(\mathbf{c})\}$.

We extend the measurability by defining that a quantity $x \in \mathbb{N}$ is *additively measurable* via sequence σ iff the sum of the contents in $\delta(\mathbf{0}, \sigma)$ is equal to x , i.e., $\sum_{i \in [n]} \delta_i(\mathbf{0}, \sigma) = x$. The set of quantities that are *additively measurable* using the capacities in \mathbf{c} is denoted by $\mathbf{M}^+(\mathbf{c})$. Obviously, this is more general than $\mathbf{M}(\mathbf{c})$ and can measure larger quantities up to $\sum_{i=1}^n c_i$. We prove that all of the additively measurable quantities again are multiples of the greatest common divisor of the capacities, that is, $\mathbf{M}^+(\mathbf{c}) = \{m \cdot \gcd(\mathbf{c}) \mid \text{for all non-negative integer } m \leq \sum_{i=1}^n c_i / \gcd(\mathbf{c})\}$.

Each $x \in \mathbf{M}^+(\mathbf{c})$ has one (or more) vector representation $\mathbf{x} = (x_1, \dots, x_n) \in \mathbb{Z}^n$ with respect to \mathbf{c} , such that $x = \mathbf{x} \cdot \mathbf{c} = \sum_{i=1}^n x_i c_i$. Moreover, we say a vector representation \mathbf{x} is optimal if $\|\mathbf{x}\|_1$ is minimum. We denote $\mu_{\mathbf{c}}(x) = \min_{\mathbf{x} \cdot \mathbf{c} = x} \|\mathbf{x}\|_1$. We will use a representation \mathbf{x} to construct a sequence of operations achieving the quantity x , and vice versa.

3. Measurability and standard sequences

3.1. Measurability

First we show that the additively measurable quantities are multiples of $\gcd(\mathbf{c})$, which form a subgroup of $(\mathbb{Z}_q, +)$ with $\gcd(\mathbf{c})$ as a generator, where $q = \sum_{i=1}^n c_i$. It can be proved by induction on the number of jugs.

Theorem 1. $\mathbf{M}^+(\mathbf{c}) = \{m \cdot \gcd(\mathbf{c}) \mid \text{for all non-negative integer } m \leq \sum_{i=1}^n c_i / \gcd(\mathbf{c})\}$, for $\mathbf{c} \in \mathbb{Z}_+^n$.

Proof. There are two parts to be proved. First we show that any additively measurable quantity is a multiple of $\gcd(\mathbf{c})$. Secondly, we prove that every non-negative multiple of $\gcd(\mathbf{c})$, bounded by $\sum_{i=1}^n c_i$, is additively measurable.

The first part is simple. Assume that x is additively measurable with \mathbf{c} . Let (s_1, \dots, s_n) be a reachable state with $x = \sum_{i=1}^n s_i$. Note that each s_i is measurable. By Proposition 1, we know each s_i is a multiple of $\gcd(\mathbf{c})$ and thus x is a multiple of $\gcd(\mathbf{c})$.

We prove the second part by induction on n , the number of jugs. It is trivial when $n = 1$. Assume that the theorem holds up to $n - 1$. Assume that $x = m \cdot \gcd(\mathbf{c})$ and $x \leq \sum_{i=1}^n c_i$, for some $m \in \mathbb{N}$. Since $c_1 \leq c_2 \leq c_3 \leq \dots \leq c_n$, we have that $c_n \geq \gcd(c_1, c_2, c_3, \dots, c_{n-1})$. If $x \leq \sum_{i=1}^{n-1} c_i$, then let $y = x \bmod \gcd(c_1, c_2, c_3, \dots, c_{n-1})$. We know that $x - y$ is a multiple of $\gcd(c_1, c_2, c_3, \dots, c_{n-1})$ and thus a multiple of $\gcd(\mathbf{c})$. We already know x is a multiple of $\gcd(\mathbf{c})$ by assumption, and then so is y . By Proposition 1, we have $y \in \mathbf{M}(\mathbf{c})$. This implies that we can reach $(0, 0, 0, \dots, y)$ first. By induction hypothesis, $x - y \in \mathbf{M}^+(c_1, c_2, \dots, c_{n-1})$. So we can reach a state s by using the first $n - 1$ jugs, where the sum of its contents is equal to $x - y$. Together with the quantity y in jug n , we can achieve the total sum x .

If $x > \sum_{i=1}^{n-1} c_i$, then let $y = x - \sum_{i=1}^{n-1} c_i \leq c_n$. We know that y is a multiple of $\gcd(\mathbf{c})$, since x and $\sum_{i=1}^{n-1} c_i$ are both multiples of $\gcd(\mathbf{c})$, and thus by Proposition 1 we have $y \in \mathbf{M}(\mathbf{c})$. This implies that we can reach $(0, 0, 0, \dots, x - \sum_{i=1}^{n-1} c_i)$ first, and then we can reach a state s with the sum of the quantity in each jug equal to x , by filling all of the jugs other than jug n . From the above, we know that $\mathbf{M}^+(\mathbf{c}) = [0, \sum_{i=1}^n c_i] \cap \{m \cdot \gcd(\mathbf{c}) \mid \forall m \in \mathbb{N}\}$. \square

3.2. Standard sequences

A sequence of operations is called a *standard* sequence if *fill* operations are always applied to empty jugs, and *empty* operations to full jugs. Formally, a standard sequence $\sigma = o_1 \dots o_m$ satisfies:

$$\forall o_l \in \{\downarrow i, i \uparrow\}, \delta_i(\mathbf{0}, o_1 \dots o_{l-1}) = \begin{cases} 0, & \text{if } o_l = \downarrow i \\ c_i, & \text{if } o_l = i \uparrow. \end{cases}$$

The total amount of water in the jugs can be changed only by non-pour operations. Since every fill operation is applied to an empty jug and every empty operation is applied to a full jug in a standard sequence, each step changing the total amount of water in the sequence either increases or decreases c_i liters for some i . Therefore, if σ is standard, then $(f_1(\sigma) - e_1(\sigma), \dots, f_n(\sigma) - e_n(\sigma))$ is a vector representation of $\sum_{i=1}^n \delta_i(\mathbf{0}, \sigma)$.

In this subsection, we show that for every reachable state s , there exists an optimal standard sequence σ such that $\delta(\mathbf{0}, \sigma) = s$. According to this fact, we prove that a sequence that additively measures x has at least $\mu_c(x)$ non-pour operations. We also give an algorithm outputs a sequence that additively measures x with exactly $\mu_c(x)$ non-pour operations. The bounds for the number of pour operations will be discussed in Section 4.

The existence of optimal standard sequences is crucial. We show that any sequence of operations can be transformed into standard one without increasing length.

Lemma 2. *Let $\sigma = o_1 o_2 \dots o_m \in O^*$ be an arbitrary sequence of m operations such that $\delta(\mathbf{0}, \sigma) = \mathbf{t} = (t_1, t_2, \dots, t_n)$. Then for every $i \in [n]$, there exists a sequence of m operations ρ such that $\delta(\mathbf{0}, \rho) \neq (t_1, \dots, t_{i-1}, 0, t_{i+1}, \dots, t_n)$ implies $\delta(\mathbf{0}, \rho) = (t_1, \dots, t_{i-1}, c_i, t_{i+1}, \dots, t_n)$.*

Proof. Let $i \in [n]$. We prove the lemma by induction on m . The base case $m = 0$ is straightforward. Assume that the lemma holds for $m < k$ and let o_l be the last operation involving jug i . Let $\mathbf{t}^l = (t_1^l, \dots, t_n^l) = \delta(\mathbf{0}, o_1 \dots o_l)$. If $t_i^l = 0$ or $t_i^l = c_i$ then the lemma holds. Otherwise, o_l must be $i \rightarrow j$ or $j \rightarrow i$ for some j . Moreover, after applying $o_1 \dots o_l$, $t_j^l = c_j$ if $o_l = i \rightarrow j$, else 0. Now let

$$q_l = \begin{cases} \downarrow j, & \text{if } o_l = i \rightarrow j \\ j \uparrow, & \text{if } o_l = j \rightarrow i. \end{cases}$$

Let $\mathbf{u} = (u_1, \dots, u_n) = \delta(\mathbf{0}, o_1 \dots o_{l-1})$. By the induction hypothesis, there exists a sequence $q_1 \dots q_{l-1}$ such that if $\delta(\mathbf{0}, q_1 \dots q_{l-1}) \neq (u_1, \dots, u_{i-1}, 0, u_{i+1}, \dots, u_n)$ then $\delta(\mathbf{0}, q_1 \dots q_{l-1}) = (u_1, \dots, u_{i-1}, c_i, u_{i+1}, \dots, u_n)$. Observe that o_l only affects jugs i and j . Therefore $u_r = t_r^l$ for every $r \in [n] - \{i, j\}$. We have that if $\delta(\mathbf{0}, q_1 \dots q_l) \neq (t_1^l, \dots, t_{i-1}^l, 0, t_{i+1}^l, \dots, t_n^l)$ then $\delta(\mathbf{0}, q_1 \dots q_l) = (t_1^l, \dots, t_{i-1}^l, c_i, t_{i+1}^l, \dots, t_n^l)$. Since o_{l+1}, \dots, o_k do not involve jug i , we have that if $\delta(\mathbf{0}, q_1 \dots q_l o_{l+1} \dots o_k) \neq (t_1, \dots, t_{i-1}, 0, t_{i+1}, \dots, t_n)$ then $\delta(\mathbf{0}, q_1 \dots q_l o_{l+1} \dots o_k) = (t_1, \dots, t_{i-1}, c_i, t_{i+1}, \dots, t_n)$. \square

The following lemma states that for every sequence σ' , we can construct a standard sequence σ which reaches the same state as σ' does without length overhead.

Lemma 3. *For any reachable state $s = (s_1, \dots, s_n)$, if $\delta(\mathbf{0}, \sigma') = \delta(\mathbf{0}, o'_1 \dots o'_m) = s$, then there exists a standard sequence $\sigma = o_1 \dots o_k$ with $k \leq m$ and $\delta(\mathbf{0}, \sigma) = s$.*

Proof. We prove the lemma by induction on $|\sigma'|$. The base case $|\sigma'| = 0$ is trivial. Assume that the lemma is true for $|\sigma'| < r$. Consider $\sigma' = o'_1 \dots o'_r$. If o'_r is a pour operation, then by the induction hypothesis, we are done. Thus we can assume that o'_r is a non-pour operation. Assume that o'_r is applied to jug i and $\mathbf{t} = (t_1, \dots, t_n) = \delta(\mathbf{0}, o'_1 \dots o'_r)$. By lemma 2, there is a sequence $\rho' = p'_1 \dots p'_{r-1}$ such that $\delta(\mathbf{0}, \rho') = (t_1, \dots, t_{i-1}, t'_i, t_{i+1}, \dots, t_n)$ where $t'_i \neq 0$ implies $t'_i = c_i$. By the induction hypothesis, there exists a standard sequence ρ , where $|\rho| \leq r - 1$ and $\delta(\mathbf{0}, \rho) = (t_1, \dots, t_{i-1}, t'_i, t_{i+1}, \dots, t_n)$.

For the case $o'_r = \downarrow i$: if $s_i = 0$ then $\rho o'_r$ is a standard sequence and $\delta(\mathbf{0}, \rho o'_r) = s$; else if $s_i \neq 0$ then $s_i = c_i$ and it means that o'_r is a redundant operation and hence ρ satisfies the lemma. Similarly, for the case $o'_r = i \uparrow$: if $s_i = 0$ then ρ satisfies the lemma; otherwise $\rho o'_r$ does. \square

Lemma 3 implies that for every reachable state s , there exists an optimal standard sequence of operations for it. Now we can derive the lower bound of non-pour operations.

Lemma 4. *For any reachable state s with $x = \sum_{i=1}^n s_i$, there exists an optimal sequence of operations $\sigma = o_1 \dots o_m$ with $\delta(\mathbf{0}, \sigma) = s$, such that the number of non-pour operations in σ is at least $\mu_c(x)$.*

Proof. By Lemma 3, there exists an optimal standard sequence of operations σ reaching s . Since σ is standard, the total amount of water in the jugs is increased or decreased by the amount c_i after $\downarrow i$ or $i \uparrow$ operation for

Algorithm MEASURE($\mathbf{c}, x, \mathbf{x}$)

Input: $\mathbf{c} = (c_1, \dots, c_n)$, the capacity of jugs.

x , the quantity to be measured.

$\mathbf{x} = (x_1, \dots, x_n)$, the optimal representation of x that achieves $\mu_{\mathbf{c}}(x)$.

Output: a sequence σ , such that $\delta^*(\mathbf{0}, \sigma)$ achieves the quantity x .

Variable: s , the state of jugs, which is initialized to be zero state.

$\mathbf{v} = (v_1, \dots, v_n)$, initialized to be \mathbf{x} .

begin

1. $\sigma := \epsilon$;

2. **for all** i **if** ($s_i = 0$ and $v_i > 0$) **do** $fill(i)$;

3. **while**($\exists i$ s.t. $v_i < 0$) **do**

4. Find j s.t. $s_j > 0$ and $v_j \geq 0$;

5. $pour(j, i)$;

6. **if** ($s_j = 0$ and $v_j > 0$) **then** $fill(j)$;

7. **if** ($s_i = c_i$) **then** $empty(i)$;

8. **while**($\exists v_i > 0$) **do**

9. Find $j > n - l$ with $v_j = 0$ and $s_j \neq c_j$;

10. $pour(i, j)$;

11. **if** $s_i = 0$ **then** $fill(i)$;

end

Procedure $fill(i)$

begin $\sigma := \sigma \circ (\downarrow i)$; $s_i := c_i$; $v_i := v_i - 1$; **end**

Procedure $empty(i)$

begin $\sigma := \sigma \circ (i \uparrow)$; $s_i := 0$; $v_i := v_i + 1$; **end**

Procedure $pour(i, j)$

begin

1. $\sigma := \sigma \circ (i \rightarrow j)$;

2. **if** ($s_i + s_j > c_j$) **then** $\{s_i := s_i + s_j - c_j; s_j := c_j\}$

3. **else** $\{s_j := s_i + s_j; s_i := 0\}$

end

Fig. 1. Measuring algorithm given $\mu_{\mathbf{c}}(x)$.

each $i \in [n]$. Note that pour operations do not change the sum. There are $f_i(\sigma)$ ($\downarrow i$)-operations and $e_i(\sigma)$ ($i \uparrow$)-operations. We have $x = \sum_{i=1}^n (f_i(\sigma) - e_i(\sigma))c_i$ and thus $\sum_{i=1}^n |f_i(\sigma) - e_i(\sigma)| \geq \mu_{\mathbf{c}}(x)$. From the above we have $\sum_{i=1}^n (f_i(\sigma) + e_i(\sigma)) \geq \sum_{i=1}^n |f_i(\sigma) - e_i(\sigma)| \geq \mu_{\mathbf{c}}(x)$. \square

Now we introduce the basic idea of our algorithm that achieves the lower bound. Assume that there is an extra jug with infinite capacity and $x = \sum_{i=1}^n c_i x_i$. Then the following operations can measure x : (1) for each $x_i > 0$ repeat $\{\downarrow i; i \rightarrow (n+1)\}$ for x_i times; (2) for each $x_i < 0$ repeat $\{(n+1) \rightarrow i; i \uparrow\}$ for $|x_i|$ times. The total number of non-pour operations is exactly $\sum_{i=1}^n |x_i|$ and the total number of measuring steps is $2 \sum_{i=1}^n |x_i|$ steps. With the above observation, given the optimal representation of x , we obtain an algorithm as in Fig. 1, which additively measures x by $\mu_{\mathbf{c}}(x)$ non-pour operations and at most $\mu_{\mathbf{c}}(x) + l - 1$ pour operations, where l is the minimum number of jugs needed to hold the quantity x . The key idea is simply simulate the imaginary jug of infinite capacity with the n jugs. However, it is not clear how to compute $\mu_{\mathbf{c}}(x)$ and the optimal representation efficiently.

Before we prove the correctness of the algorithm, we give a small execution example of the algorithm. Assume that $\mathbf{c} = (3, 15, 16)$ and we want to additively measure $x = 5$ corresponding to the representation $(2, 1, -1)$. The

algorithm MEASURE(c, x, \mathbf{x}) outputs a sequence $\sigma = \downarrow 1 \circ \downarrow 2 \circ 2 \rightarrow 3 \circ 1 \rightarrow 3 \circ 3 \uparrow \circ 1 \rightarrow 3 \circ \downarrow 1$. The state transitions of σ are shown in the following table.

1st loop	$(0, 0, 0) \xrightarrow{\downarrow 1} (3, 0, 0) \xrightarrow{\downarrow 2} (3, 15, 0)$
2nd loop	$(3, 15, 0) \xrightarrow{2 \rightarrow 3} (3, 0, 15) \xrightarrow{1 \rightarrow 3} (2, 0, 16) \xrightarrow{3 \uparrow} (2, 0, 0)$
3rd loop	$(2, 0, 0) \xrightarrow{1 \rightarrow 3} (0, 0, 2) \xrightarrow{\downarrow 1} (3, 0, 2)$

We prove the correctness in this section and leave the analysis on the number of operations in the output sequence in Section 4. We prove the correctness with the following lemma.

Lemma 5. *Let \mathbf{x} be an optimal representation of x with capacity c . The algorithm MEASURE outputs a sequence of operations σ such that $\delta(\mathbf{0}, \sigma) = s$ and $\sum_{i=1}^n s_i = x$.*

Proof. Let $F = \{i | x_i > 0\}$ and $E = \{i | x_i < 0\}$. The integer variables v_i 's are used to track the number of empty and fill operations performed. Initially, $v_i = x_i$ for $i \in [n]$. After each fill operation some v_i with $i \in F$ will decrease by 1, and after each empty operation some v_j with $j \in E$ will increase by 1. Observe that during the execution the number of fill and empty operations performed on jug i is $(x_i - v_i)$ for $i \in F$ and $(v_i - x_i)$ for $i \in E$, respectively. Thus the total quantity in the jugs is $\sum_{i \in F} (x_i - v_i)c_i - \sum_{i \in E} (v_i - x_i)c_i$ during the operations.

After the first loop (in line 2), $s_i = c_i$ for all $i \in F$. Next we show that after an iteration of the second loop (in lines 3–7), if there still exists a $v_i < 0$, then we can always find j in line 4. There are two possibilities after $\text{pour}(j, i)$ is executed in line 5, i.e., jug j can become non-empty or empty.

- (Case 1: Jug j is still non-empty.) If the loop condition holds, we can always find j in line 4, since $s_j > 0$ and $v_j \geq 0$.
- (Case 2: Jug j becomes empty.) If $v_j > 0$, then jug j will be refilled immediately and $s_j > 0$. If $v_j = 0$, suppose that line 4 fails to find a j in the subsequent iteration, then it implies that for all $v_k \geq 0, s_k = 0$. However, there does not exist such $v_k > 0$ since jug k would be refilled right after line 6. Thus v_k must be 0 for all $k \in F$. Line 7 shows that for all $i \in E$, if $v_i < 0$ then $s_i < c_i$ and thus the amount of water in the jugs is less than $\sum_{i \in E; v_i < 0} c_i \cdot c_i$. Since we have done $\sum_{i \in F} x_i$ fill operations and $\sum_{i \in E} (v_i - x_i)$ empty operations, the quantity of water left in the jugs is exactly $\sum_{i \in F} c_i x_i + \sum_{i \in E} c_i (x_i - v_i) = \mathbf{x} \cdot \mathbf{c} - \sum_{i \in E} c_i v_i$, which is greater than $\sum_{i \in E; v_i < 0} c_i - a$ contradiction!

Thus, as long as the loop condition in line 3 holds, one can always find (in line 4) a jug to perform the pour operation.

After the second loop (lines 3–7), no more empty operation will be performed. For all $i \in E, v_i = 0$ and for all $i \in F$ with $v_i > 0$, we have $s_i > 0$ by line 6. Note that the quantity of water left in the jugs is $\sum_{i \in F} (x_i - v_i)c_i - \sum_{j \in E} (v_j - x_j)c_j = x - \sum_{i \in F} c_i v_i > 0$. If for all $i \in F, v_i = 0$, then we are done.

By the assumption, the largest l jugs are sufficient to contain the quantity x . Thus, $\sum_{j > n-l} c_j \geq x$. Note that we always fill jug i at line 6 and at line 11 when $v_i > 0$ and $s_i = 0$. We have that for every $k \in [n]$ if $v_k > 0$ then $s_k > 0$ at line 8. This implies that there must be some jug $j > n - l$ with $s_j \neq c_j$ and $v_j = 0$. Otherwise $\forall j > n - l$ we have $s_j = c_j$ or $v_j > 0$, which implies the quantity $x = \sum_{j=1}^n (s_j + v_j c_j) \geq s_i + \sum_{j > n-l} c_j > x$ (note that $v_i > 0$ and $s_i > 0$), a contradiction. Thus, whenever the loop condition at line 8 holds (i.e., $v_i > 0$), one can always find a suitable jug for pouring at line 9.

Finally, when the algorithm terminates, it actually performed $\sum_{i \in F} x_i$ fill operations and $-\sum_{j \in E} x_j$ empty operations and the net quantity is $\sum_{i \in F} c_i x_i + \sum_{j \in E} c_j x_j = x$. \square

We remark that the output sequence of the MEASURE algorithm is a standard sequence achieving the quantity x . It matches the lower bound of non-pour operations when we use the optimal representation as input.

4. Lower and upper bounds of measuring steps

4.1. Bounds for additively measurable quantities

In this subsection, we give the lower and upper bounds for additively measuring. First, we prove the following theorem as a consequence of the lower bound on the number of pour operations in an optimal standard sequence.

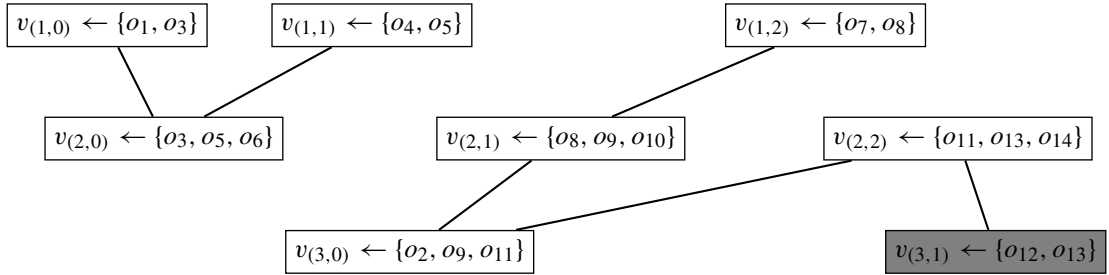


Fig. 2. \$G_{\sigma'}\$: the corresponding graph of \$\sigma'\$.

Theorem 6. Let \$s = (s_1, \dots, s_n)\$ be a reachable state, \$x = \sum_{i=1}^n s_i\$ and \$n_{ne}\$ be the number of non-zero entries of \$s\$, then no sequence of operations can reach \$s\$ in less than \$2\mu_c(x) - n_{ne}\$ steps.

We prove the lower bound for pour operations by inspecting a graph \$G_{\sigma}\$ corresponding to a standard sequence of operations \$\sigma = o_1 \dots o_m\$. Recall that \$\delta_i(s, \rho)\$ is the quantity in jug \$i\$ after applying a sequence of operations \$\rho\$ to state \$s\$. We denote the number of operations making jug \$i\$ empty in \$\sigma\$ as \$z_i(\sigma)\$, i.e., \$z_i(\sigma) = |\{k : k > 0, \delta_i(\mathbf{0}, o_1 \dots o_k) = 0, \text{ and } \delta_i(\mathbf{0}, o_1 \dots o_{k-1}) > 0\}|\$. Note that \$z_i(\sigma)\$ can be larger than \$e_i(\sigma)\$.

For jug \$i\$, we define vertices \$v_{(i,j)}\$, where \$i \in [n]\$ and \$j = 0, \dots, z_i(\sigma)\$ if \$\delta_i(\mathbf{0}, \sigma) \neq 0\$, else \$j = 0, \dots, (z_i(\sigma) - 1)\$. Let \$V_{\sigma}\$ be the set of all \$v_{(i,j)}\$'s. For each vertex \$v_{(i,j)} \in V_{\sigma}\$, we associate it with every operation \$o_k\$ satisfying: (1) \$o_k\$ is applied to jug \$i\$. (2) There are exactly \$j\$ operations making jug \$i\$ empty before \$o_k\$. More precisely, we associate \$v_{(i,j)}\$ with the following set of operations:

$$\{o_k : o_k \text{ involves jug } i \text{ and } |\{o_{k'} : k' < k, \delta_i(\mathbf{0}, o_1 \dots o_{k'}) = 0 \text{ and } \delta_i(\mathbf{0}, o_1 \dots o_{k'-1}) > 0\}| = j\}.$$

Note that in a standard sequence each \$v_{(i,j)}\$ is associated with at most one empty and one fill operation. Now we define the corresponding graph \$G_{\sigma} = \langle V_{\sigma}, E_{\sigma} \rangle\$ of \$\sigma\$, where \$V_{\sigma}\$ is the set of all \$v_{(i,j)}\$'s and \$\{v_{(i,j)}, v_{(i',j')}\} \in E_{\sigma}\$ if and only if both \$v_{(i,j)}\$ and \$v_{(i',j')}\$ are associated with a common operation \$o_l\$. Therefore, \$|E_{\sigma}|\$ is the number of pour operations. If every operation in the set associated with \$v_{(i,j)}\$ does not make jug \$i\$ empty, then we color \$v_{(i,j)}\$ gray, otherwise we color it white.

There are some interesting properties about the first and last operations associated with a vertex. Let \$o_s\$ and \$o_e\$ be the first and last operations associated with vertex \$v_{(i,j)} \in V_{\sigma}\$. We have \$\delta_i(\mathbf{0}, o_1 \dots o_{s-1}) = 0\$, since if there exists an operation applied to jug \$i\$ before \$o_s\$, then it must make jug \$i\$ empty (note that if \$o_s\$ is the first operation applied to jug \$i\$ then prior to \$o_s\$ jug \$i\$ is empty). Also \$\delta_i(\mathbf{0}, o_1 \dots o_e) = 0\$ if and only if \$v_{(i,j)}\$ is white (note that a gray vertex cannot be associated with an operation making it empty).

For example, we consider an instance with \$c = (14, 28, 31)\$ and \$x = 20\$. Let \$\sigma' = o_1 o_2 \dots o_{14} = \downarrow 1 o \downarrow 3 o 1 \rightarrow 2 o \downarrow 1 o 1 \rightarrow 2 o 2 \uparrow o \downarrow 1 o 1 \rightarrow 2 o 3 \rightarrow 2 o 2 \uparrow o 3 \rightarrow 2 o \downarrow 3 o 3 \rightarrow 2 o 2 \uparrow\$. It can be justified that \$\delta(\mathbf{0}, \sigma') = (0, 0, 20)\$, but \$\sigma'\$ is not an optimal sequence. The detailed transitions by \$\sigma'\$ are: \$(0, 0, 0) \xrightarrow{\downarrow 1} (14, 0, 0) \xrightarrow{\downarrow 3} (14, 0, 31) \xrightarrow{1 \rightarrow 2} (0, 14, 31) \xrightarrow{\downarrow 1} (14, 14, 31) \xrightarrow{1 \rightarrow 2} (0, 28, 31) \xrightarrow{2 \uparrow} (0, 0, 31) \xrightarrow{\downarrow 1} (14, 0, 31) \xrightarrow{1 \rightarrow 2} (0, 14, 31) \xrightarrow{3 \rightarrow 2} (0, 28, 17) \xrightarrow{2 \uparrow} (0, 0, 17) \xrightarrow{3 \rightarrow 2} (0, 17, 0) \xrightarrow{\downarrow 3} (0, 17, 31) \xrightarrow{3 \rightarrow 2} (0, 28, 20) \xrightarrow{2 \uparrow} (0, 0, 20)\$. We construct the corresponding graph in Fig. 2, where each vertex \$v\$ is associated with a set of operations \$O_v\$ (denoted as \$v \leftarrow O_v\$ in a rectangle) and \$z_1(\sigma') = z_2(\sigma') = 3, z_3(\sigma') = 1\$.

We now prove the following crucial lemma, which is a key tool to prove the lower bound of pour operations.

Lemma 7. Let \$G_{\sigma} = \langle V_{\sigma}, E_{\sigma} \rangle\$ be the corresponding graph of \$\sigma = o_1 \dots o_m\$. If a connected component \$G'_{\sigma} = \langle V'_{\sigma}, E'_{\sigma} \rangle\$ of \$G_{\sigma}\$ contains no gray vertex, then \$\delta(\mathbf{0}, \sigma) = \delta(\mathbf{0}, \sigma')\$, where \$\sigma'\$ is obtained from \$\sigma\$ by removing all operations in the sets associated with the vertices in \$V'_{\sigma}\$.

Proof. Suppose that \$\sigma' = o'_1 \dots o'_{m'}\$. Since \$\sigma'\$ is obtained by removing some operations from \$\sigma\$, we define a one-to-one mapping \$g : \{1, \dots, m'\} \to \{1, \dots, m\}\$ such that \$o'_i\$ in \$\sigma'\$ corresponds to \$o_{g(i)}\$ in \$\sigma\$, for \$i = 1\$ to \$m'\$. In other words, \$\sigma'\$ is simply a projection of \$\sigma\$. We claim that: for every \$k \in \{1, \dots, m'\}\$, if \$o'_k\$ is applied to jug \$i\$, then \$\delta_i(\mathbf{0}, o'_1 \dots o'_k) = \delta_i(\mathbf{0}, o_1 \dots o_{g(k)})\$.

We show that $\delta(\mathbf{0}, \sigma) = \delta(\mathbf{0}, \sigma')$ follows from the claim. First consider every i with $\delta_i(\mathbf{0}, \sigma') \neq 0$. Let o'_l be the last operation applied to jug i in σ' . Since $\delta_i(\mathbf{0}, \sigma') \neq 0$, $v_{(i, z_i(\sigma))}$ is gray. The last operation applied to jug i is still in σ' . Therefore, there is no operation o_{l^*} applied to jug i in σ with $l^* > g(l)$. We have that $\delta_i(\mathbf{0}, \sigma') = \delta_i(\mathbf{0}, o'_1 \cdots o'_l) = \delta_i(\mathbf{0}, o_1 \cdots o_{g(l)}) = \delta_i(\mathbf{0}, \sigma)$. Next for those i with $\delta_i(\mathbf{0}, \sigma) = 0$, there are two possibilities: (1) All operations on jug i are removed. Therefore, $\delta_i(\mathbf{0}, \sigma') = \delta_i(\mathbf{0}, \epsilon) = 0 = \delta_i(\mathbf{0}, \sigma)$; (2) There are still some operations on jug i in σ' . Note that $\delta_i(\mathbf{0}, \sigma) = 0$ implies that all operations on jug i in σ are associated with white vertices. Let the last operation applied to jug i in σ' be o'_l . Then $\delta_i(\mathbf{0}, \sigma') = \delta_i(\mathbf{0}, o'_1 \cdots o'_l) = \delta_i(\mathbf{0}, o_1 \cdots o_{g(l)}) = 0 = \delta_i(\mathbf{0}, \sigma)$, since $o_{g(l)}$ is the last operation associated with some white vertex. We conclude that for every $i \in [n]$, $\delta_i(\mathbf{0}, \sigma') = \delta_i(\mathbf{0}, \sigma)$ which implies $\delta(\mathbf{0}, \sigma') = \delta(\mathbf{0}, \sigma)$.

Now we prove the claim by the way of contradiction. Let k be the smallest index such that when o'_k is applied to some jug π , $\delta_\pi(\mathbf{0}, o'_1 \cdots o'_k) \neq \delta_\pi(\mathbf{0}, o_1 \cdots o_{g(k)})$. We have that o'_k must be a pour operation, thus we assume that $o'_k = i \rightarrow j$. Since o'_k is applied to both jug i and jug j , one of the following two inequalities must hold: $\delta_i(\mathbf{0}, o'_1 \cdots o'_{k-1}) \neq \delta_i(\mathbf{0}, o_1 \cdots o_{g(k-1)})$ or $\delta_j(\mathbf{0}, o'_1 \cdots o'_{k-1}) \neq \delta_j(\mathbf{0}, o_1 \cdots o_{g(k-1)})$. Therefore, if we can show that $\delta_i(\mathbf{0}, o'_1 \cdots o'_{k-1}) = \delta_i(\mathbf{0}, o_1 \cdots o_{g(k-1)})$ and $\delta_j(\mathbf{0}, o'_1 \cdots o'_{k-1}) = \delta_j(\mathbf{0}, o_1 \cdots o_{g(k-1)})$, then the assumption leads to a contradiction.

We show that $\delta_i(\mathbf{0}, o'_1 \cdots o'_{k-1}) = \delta_i(\mathbf{0}, o_1 \cdots o_{g(k-1)})$, and the proof for $\delta_j(\mathbf{0}, o'_1 \cdots o'_{k-1}) = \delta_j(\mathbf{0}, o_1 \cdots o_{g(k-1)})$ is similar. If $o_{g(k)}$ is the first operation applied to jug i in σ , then $\delta_i(\mathbf{0}, o_1 \cdots o_{g(k-1)}) = \delta_i(\mathbf{0}, \epsilon) = 0 = \delta_i(\mathbf{0}, o'_1 \cdots o'_{k-1})$. Assume that $o_{g(k)}$ is not the first operation applied to jug i . Let o_l be the operation applied to jug i prior to $o_{g(k)}$. We have two possible cases:

- (Case 1: o_l is not removed.) Then we have $\delta_i(\mathbf{0}, o_1 \cdots o_{g(k-1)}) = \delta_i(\mathbf{0}, o_1 \cdots o_l) = \delta_i(\mathbf{0}, o'_1 \cdots o'_{g-1(l)}) = \delta_i(\mathbf{0}, o'_1 \cdots o'_{k-1})$, since by the assumption there is no integer $k' < k$ such that $o_{k'}$ is applied to jug i and $\delta_i(\mathbf{0}, o'_1 \cdots o'_{k'}) \neq \delta_i(\mathbf{0}, o_1 \cdots o_{g(k')})$.
- (Case 2: o_l is removed.) Then o_l must make jug i empty, otherwise $o_{g(k)}$ and o_l would be associated with the same removed vertex. If all operations involving jug i before $o_{g(k)}$ are removed, then $\delta_i(\mathbf{0}, o_1 \cdots o_{g(k-1)}) = 0 = \delta_i(\mathbf{0}, o'_1 \cdots o'_{k-1})$. Assume that some of them are not removed. Let $o_{g(h)}$ be the last remaining operation applied to jug i before $o_{g(k)}$. Since there is no other operation applied to jug i between $o_{g(h)}$ and $o_{g(k)}$, $o_{g(h)}$ is the last operation associated with some white vertex. Thus, $o_{g(h)}$ empties jug i and $\delta_i(\mathbf{0}, o_1 \cdots o_{g(k-1)}) = 0 = \delta_i(\mathbf{0}, o_1 \cdots o_{g(h)}) = \delta_i(\mathbf{0}, o'_1 \cdots o'_h) = \delta_i(\mathbf{0}, o'_1 \cdots o'_{k-1})$.

Thus $\delta_i(\mathbf{0}, o'_1 \cdots o'_{k-1}) = \delta_i(\mathbf{0}, o_1 \cdots o_{g(k-1)})$ and similarly $\delta_j(\mathbf{0}, o'_1 \cdots o'_{k-1}) = \delta_j(\mathbf{0}, o_1 \cdots o_{g(k-1)})$. Therefore the claim is true and the lemma is proved. \square

Now we can give the lower bound on the number of pour operations in a standard sequence.

Lemma 8. Let $\sigma = o_1 \cdots o_m$ be a standard sequence for a reachable state $s = \delta(\mathbf{0}, \sigma)$ and $x = \sum_{i=1}^n s_i$. Let n_{ne} be the number of non-zero entries of s , then the number of pour operations in σ is at least $\mu_c(x) - n_{ne}$.

Proof. Let V_σ^w be the vertex set of all connected components of G_σ that contains no gray vertex. By removing all operations associated with vertices in V_σ^w , we obtain a standard sequence ρ such that $\delta(\mathbf{0}, \rho) = s$, $p(\rho) \leq p(\sigma)$ and its corresponding graph $G_\rho = \langle V_\rho, E_\rho \rangle$ does not contain any connected component consisting of only white vertices. Since G_ρ has at most n_{ne} connected components, we have $p(\rho) = |E_\rho| \geq |V_\rho| - n_{ne}$. Since ρ is standard, for every vertex $v \in V_\rho$, v is associated with at most one fill operation and at most one empty operation. We have $|V_\rho| \geq \sum_{i=1}^n \max(e_i(\rho), f_i(\rho)) \geq \sum_{i=1}^n |e_i(\rho) - f_i(\rho)| \geq \mu_c(x)$; hence the number of pour operations in σ is at least $\mu_c(x) - n_{ne}$. \square

By Lemmas 4 and 8, we have Theorem 6 as an immediate consequence. We remark that this lower bound is tight empirically for many cases, for example, for measuring $(0, \dots, 0, x)$ with $x \in \mathbf{M}(\mathbf{c})$.

Now we turn to the upper bound. We show the following theorem by analyzing the MEASURE algorithm carefully.

Theorem 9. For all $x \in M^+(\mathbf{c})$, if we can use the largest l jugs to hold the quantity x , then the algorithm MEASURE additively measures x in $2\mu_c(x) + l - 1$ steps.

Since we already know the number of non-pour operations generated by the algorithm is $\|\mathbf{x}\|_1$, we focus on the number of pour operations.

Lemma 10. Let x be an optimal representation of x with capacity c . The algorithm MEASURE outputs a sequence of operations σ such that $|\sigma| \leq 2\mu_c(x) + l - 1$.

Proof. Let $F = \{i | x_i > 0\}$ and $E = \{i | x_i < 0\}$. Let e be the number of empty operation in σ . There are 3 loops in algorithm MEASURE. Let f_k be the numbers of fill operations generated in the k -th loop, $k = 1, 2, 3$. Since for every $i \in F$, we fill jug i in the first loop, we have $f_1 = |F|$.

If an invocation of $pour(j, i)$ in the second loop is neither followed by an $empty(i)$ nor by a $fill(j)$, then we have $v_j = 0$ and $s_j = 0$ at the end of that iteration. This means that we will never pick such j in line 4 afterward. Note that we only pick j 's from F . Therefore, the number of pour operations generated in the second loop is at most $f_2 + e + |F| = f_2 + e + f_1$.

If we do not fill jug i after invoking $pour(i, j)$ in the third loop, then we have jug j is full and jug i is non-empty. Such j will never be picked in line 9 again. Moreover, such event happens at most $l - 1$ times, otherwise $x > \sum_{k=n-l+1}^n c_k$, a contradiction. The total number of pour operations generated in the third loop is at most $f_3 + l - 1$. The total number of pour operations in σ is no more than $f_2 + e + f_1 + f_3 + l - 1 = \|\mathbf{x}\|_1 + l - 1 = \mu_c(x) + l - 1$ and thus $|\sigma| \leq 2\mu_c(x) + l - 1$. \square

Theorem 9 follows by Lemmas 5 and 10. The above proof is suggested by one of the anonymous referees which simplifies our proof in the early draft.

4.2. Bounds for measurable quantities

In this subsection, we prove new lower and upper bounds for measurable quantities. Recall that a sequence σ measures x if one of the entries of $\delta(\mathbf{0}, \sigma)$ equal to x , but a sequence σ' additively measures x if the sum over all entries of $\delta(\mathbf{0}, \sigma')$ equal to x . Therefore, we cannot directly apply the results for additively measuring here. We prove the following lower bound, which improves the previous bound $\frac{1}{2}\mu_c(x)$ by Boldi et al. [4].

Theorem 11. No sequence of operations can measure x , for all $x \in M(c)$, in less than $\max\{2\mu_c(x) - n, \mu_c(x)\}$ steps.

Proof. First we show that no sequence can measure x in less than $2\mu_c(x) - n$ steps by way of contradiction. If there is a state s with $s_i = x$ for some $i \in [n]$ that can be reached in less than $2\mu_c(x) - n$ steps, then $(\overbrace{0, \dots, 0}^{i-1}, x, 0, \dots, 0)$ is reachable in less than $2\mu_c(x) - 1$ steps by applying at most $n - 1$ extra empty operations. By Theorem 6, we know it is impossible!

Note that the above bound vanishes when $2\mu_c(x) - n < 0$. Next we prove that no sequence measures x in less than $\mu_c(x)$ steps by a more careful observation on the number of pour operations. It is trivial for $x = 0$. Thus we consider the case when $x > 0$. Let ρ be a shortest standard sequence that measures x . Let $\delta(\mathbf{0}, \rho) = \mathbf{t}$ and $t_k = x$ for some k . Due to Lemma 2, we can assume that $t_j \in \{0, c_j\}$ for every $j \neq k$. Let $R = \{j : t_j \neq 0 \wedge j \neq k\} = \{r_1, \dots, r_{|R|}\}$.

Let $\sigma = \rho \circ r_1 \uparrow \circ \dots \circ r_{|R|} \uparrow$. Thus σ is standard, $\delta(\mathbf{0}, \sigma) = (\overbrace{0, \dots, 0}^{k-1}, x, 0, \dots, 0)$ and $p(\rho) = p(\sigma)$. By Lemma 8, we obtain $p(\rho) = p(\sigma) \geq \mu_c(x) - 1$. Since we need at least one fill operation for measuring $x > 0$, we have $|\rho| \geq 1 + p(\rho) = \mu_c(x)$. \square

Next we prove the following upper bound which improves the previous bound $\frac{5}{2}\mu_c(x)$ by Boldi et al..

Theorem 12. For all $x \in M(c)$, there exists a sequence of operations ρ such that $\delta(\mathbf{0}, \rho) = (0, \dots, 0, x)$ and $|\rho| \leq 2\mu_c(x)$.

Proof. Let $\sigma = o_1 \cdots o_{|\sigma|}$ be the sequence output by MEASURE. A pour operation $i \rightarrow j$ in σ makes jug j full or jug i empty. Let P_1 be the set of pour operations in the former case and P_2 in the latter case.

Consider an arbitrary operation $o_k = i \rightarrow j$ in P_1 . If it is generated in the second loop, then MEASURE invokes $empty(j)$ in the same iteration. Since the total amount of water is less than x in any iteration of the third loop and the largest jug has capacity at least x , one can conclude that there is no pour operation in P_1 generated in the third loop. Thus, $|P_1| = e(\sigma)$.

Let $o_k = i \rightarrow j$ be a pour operation in P_2 . Let $l = \max\{q : q < k \wedge o_q = \downarrow i\}$. We claim that for every $r \in (l, k)$, o_r is not a pour operation emptying jug i . Suppose not, let $o_r = i \rightarrow j'$ in P_2 . If o_r is generated in the third loop, then

it must be followed by $\downarrow i$, since jug i is emptied by o_r . If o_r is generated in the second loop and o_r is not followed by $\downarrow i$, then $v_i = 0$ and $s_i = 0$ after invoking $\text{pour}(i, j')$. Consequently, o_r is the last pour operation emptying jug i and this contradicts the assumption $r < k$. Then o_r must be followed by $\downarrow i$, and it leads to $r < l$, which contradicts that $r \in (l, k)$. Now we know the claim is true, and every pour operation $i \rightarrow j$ in P_2 can be paired with the closest prior $\downarrow i$ operation. Thus $|P_2| = f(\sigma) - f_u$, where f_u is the number of unpaired fill operations.

Now we turn to bound f_u . Let $F^* = \{i : i < n \wedge s_i > 0 \text{ when MEASURE terminated}\}$. $F^* \subseteq \{i : x_i > 0\}$ since at the end of the second loop for every $i \in \{i : x_i \leq 0\}$ $s_i = 0$ and we only pour water into jug n in the third loop. For every $i \in F^*$, there exists $\downarrow i$ in σ , and the last one is unpaired, since there is no operation emptying jug i after it. We have $|F^*| \leq f_u$, so $|\sigma| = f(\sigma) + e(\sigma) + p(\sigma) \leq f(\sigma) + e(\sigma) + |P_1| + |P_2| \leq 2(f(\sigma) + e(\sigma)) - |F^*| \leq 2\mu_c(x) - |F^*|$. Let $\rho = \sigma \circ i_1 \rightarrow n \circ \dots \circ i_{|F^*|} \rightarrow n$ where $F^* = \{i_1, \dots, i_{|F^*|}\}$. We have that $\delta(\mathbf{0}, \rho) = (0, \dots, 0, x)$ and $|\rho| \leq 2\mu_c(x)$. \square

We thank one of the anonymous referees whose comments inspired us to give stronger results for measurable quantities than the original version. From Theorem 12, we know that when measuring $(0, \dots, 0, x)$ with $x \in \mathbf{M}(c)$, the bound is very close to the lower bound mentioned in Section 4.1.

5. Computing $\mu_c(x)$ and its approximation

5.1. Hardness of jug measuring

Recall that we define the optimal jug measuring problem as: given n jugs with capacity c_1, \dots, c_n , and a state s , find the length of shortest sequence σ such that $\delta(\mathbf{0}, \sigma) = s$.

We have used $\mu_c(x)$ to bound the number of steps on jug measuring. In this section, we investigate the hardness of computing $\mu_c(x)$ and the optimal jug measuring problem. It can be shown directly that computing $\mu_c(x)$ is indeed NP -hard. For the notations of computational complexity we refer to standard textbooks such as by Sipser [10] and Papadimitriou [9]. However, we have stronger results, i.e., it is hard even to approximate it.

To study the complexity of computing $\mu_c(x)$, we investigate the shortest GCD multiplier problem [6], which is: given c_1, c_2, \dots, c_n , we want to find $\mathbf{x} = (x_1, x_2, \dots, x_n)$ such that $\sum_{i=1}^n x_i c_i = \text{gcd}(c_1, c_2, \dots, c_n)$ and $\|\mathbf{x}\|_p$ is minimal. The latter problem, for $p = 1$, is a special case of the problem of computing $\mu_c(x)$. Havas and Seifert [6] proved that: (1) Unless $NP \subseteq P$, there exists no polynomial-time algorithm which approximates the shortest GCD multiplier problem in l_p -norm within a factor of k , where $k \geq 1$ is an arbitrary constant. (2) Unless $NP \subseteq DTIME(n^{\text{poly}(\log n)})$, there exists no polynomial-time algorithm which approximates the shortest GCD multiplier problem in l_p -norm within a factor of $n^{1/(p \log^\gamma n)}$, where γ is an arbitrary small positive constant.

Let σ be an optimal sequence to measure $x \in \mathbf{M}(c)$. By the lower and upper bounds of the previous section, we have $\mu_c(x) = \lceil \frac{|\sigma|}{2} \rceil$ where $\delta(\mathbf{0}, \sigma) = (0, \dots, 0, x)$. In other words, if we know how to solve the optimal jug measuring problem, then we know the value of $\mu_c(x)$. By selecting $x = \text{gcd}(c)$, for the optimal jug measuring problem we have analogous results:

Theorem 13. *Unless $NP \subseteq P$, there exists no polynomial-time algorithm which approximates the optimal jug measuring problem within a factor of k , where $k \geq 1$ is an arbitrary constant.*

Theorem 14. *Unless $NP \subseteq DTIME(n^{\text{poly}(\log n)})$, there exists no polynomial-time algorithm which approximates the optimal jug measuring problem within a factor of $n^{1/\log^\gamma n}$, where γ is an arbitrary small positive constant.*

5.2. Reduction from computing $\mu_c(x)$ to CVP

In this section we give a polynomial-time reduction from the problem of computing $\mu_c(x)$ to CVP and an LLL-based approximation algorithm for computing $\mu_c(x)$ with exponential errors (approximation ratio). Our approximation algorithm is based on the fact: computing $\mu_c(x)$ can be polynomially reduced to the closest lattice vector problem (CVP). This is an extension of the approach given by Havas et al. [7] for approximating the extended GCD problem.

We introduce lattice and the closest lattice problem briefly as follows. A lattice in \mathbb{R}^n is the set of all integer linear combination of m independent column vectors $\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_m$. The lattice generated by $\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_m$, denoted as

$L(\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_m)$, is the set $\{\sum_{i=1}^m \lambda_i \mathbf{b}_i \mid \forall i \in [m], \lambda_i \in \mathbb{Z}\}$. The independent vectors $\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_m$ are called a basis of the lattice. The *closest lattice vector problem* is: given a basis $\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_m$, a vector $v \in \mathbb{R}^n$ and an integer p , we want to find the lattice point $u \in L(\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_m)$ which is closest to v under l_p -norm.

In order to complete the reduction from computing $\mu_c(x)$ to CVP, we introduce the *Hermite normal form* [11]. A matrix A is said to be in Hermite normal form if it has the form $[B \ 0]$ where the matrix B is a non-singular, lower triangular, non-negative matrix, in which each row has a unique maximum entry, which is located on the main diagonal of B .

The following operations on a matrix are called *elementary (unimodular) column operations*: (1) exchanging two columns; (2) multiplying a column by -1 ; (3) adding an integral multiple of one column to another column.

A non-singular matrix U is a unimodular matrix if U is integral and has determinant 1 or -1 . Unimodular matrix

can be obtained by applying some unimodular operations to I . For example, $\begin{bmatrix} 3 & 1 & 0 \\ -1 & 0 & 0 \\ 0 & 0 & 1 \end{bmatrix}$ is a unimodular matrix

which is obtained by: exchanging column 1 and 2, multiplying column 1 by -1 , then adding 3 times column 2 to column 1. There are several known facts about Hermite normal form:

Theorem 15 ([11]). (1) Each rational matrix of full row rank can be brought into Hermite normal form by a series of elementary column operations. (2) For each rational matrix A of full row rank, there is a unimodular matrix U such that AU is the Hermite normal form of A . (3) Given a feasible system $A\mathbf{y} = \mathbf{b}$ of rational linear diophantine equations, we can find in polynomial-time integral vectors $y_0, y_1, y_2, \dots, y_t$ such that $\{\mathbf{y} \mid A\mathbf{y} = \mathbf{b}; \mathbf{y} \text{ is integral}\} = \{y_0 + \lambda_1 y_1 + \dots + \lambda_t y_t \mid \lambda_1, \dots, \lambda_t \in \mathbb{Z}\}$ with y_1, y_2, \dots, y_t linearly independent. Moreover, $[y_0 \ y_1 \ y_2 \ \dots \ y_t] = U \begin{bmatrix} B^{-1}\mathbf{b} & 0 \\ 0 & I \end{bmatrix}$, where $AU = [B \ 0]$ is the Hermite normal form of A .

We are now in position to prove the announced reduction.

Corollary 16. The problem of computing $\mu_c(x)$ can be polynomially reduced to CVP.

Proof. Assume that $((c_1, c_2, \dots, c_n), x)$ is an instance for computing $\mu_c(x)$. Let 1-by- n -matrix $C = [c_1 \ c_2 \ \dots \ c_n]$. By Theorem 15(2), there exists a unimodular n -by- n -matrix U such that $CU = [b \ 0 \ \dots \ 0]$ is the Hermite normal form of C . Let $[v_0 \ v_1 \ v_2 \ \dots \ v_{n-1}] = U \begin{bmatrix} \frac{x}{b} & 0 \\ 0 & I \end{bmatrix}$. By Theorem 15(3), we know that v_1, v_2, \dots, v_{n-1} are linearly independent column vectors and form a basis of a lattice L . Hence $\mu_c(x) = \min_{v \in \{x \mid c \cdot x = x\}} \|v\|_1 = \min_{\lambda_1, \dots, \lambda_{n-1} \in \mathbb{Z}} \|v_0 + \lambda_1 v_1 + \dots + \lambda_{n-1} v_{n-1}\|_1 = \min_{w \in L} \|w - (-v_0)\|_1$. Thus $\mu_c(x)$ is the l_1 -norm of the vector $v \in L$ which is closest to $-v_0$. Note that all computation can be done in polynomial time. \square

For the problem of computing the unimodular matrix U such that $[c_1 \ c_2 \ \dots \ c_n]U$ is in the Hermite normal form, we propose an algorithm which is simpler than the general algorithm in [11] and the algorithms mentioned in [7]. However, the algorithms in [7] could outperform our algorithm. Our algorithm works for our application on the special case $[c_1 \ c_2 \ \dots \ c_n]U$, but it cannot compute the unimodular matrix U for any other n by m matrix, where $n > 1$. The algorithm is shown in Fig. 3, and it is based on the Euclidean algorithm. It computes the greatest common divisor of the first and the i th entries by applying Euclidean algorithm with unimodular operations. Each iteration terminates when the greatest common divisor is written back to the first entry and 0 is written to the i th entry, thus it runs in $O(n^2 \log c_n)$ time. The maximum of the absolute value of the entries of U will not exceed $O(n \log c_n)$.

Babai [2] provided two polynomial-time approximation algorithms for CVP. Both algorithms are based on LLL basis reduction algorithm proposed by A. K. Lenstra, H. W. Lenstra and L. Lovasz (see e.g., [8]). Assume that we are given an LLL reduced basis $\{\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_n\}$ where $\|\mathbf{b}_1\|_2 \leq 2^{\frac{n-1}{2}} \min_{v \in L(\mathbf{b}_1, \dots, \mathbf{b}_n) - \{0\}} \|v\|_2$, a vector $\mathbf{x} = \sum_{i=1}^n \alpha_i \mathbf{b}_i$ and we are to find a vector $\mathbf{w} \in L(\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_n)$ that is close to \mathbf{x} . The first algorithm is called the rounding off heuristic algorithm. It simply outputs $\mathbf{w} = \sum_{i=1}^n \beta_i \mathbf{b}_i$ where β_i is the closest integer to α_i . The second algorithm is called the nearest-plane heuristic algorithm. It is a recursive algorithm. Let $V = \text{span}(\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_{n-1})$, and find $\mathbf{v} \in L(\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_n)$ such that the distance between $V + \mathbf{v}$ and \mathbf{x} is minimal. Let \mathbf{x}' be the orthogonal projection of \mathbf{x} on $V + \mathbf{v}$. Then find $\mathbf{y} \in L(\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_{n-1})$ close to $\mathbf{x}' - \mathbf{v}$, and output $\mathbf{w} = \mathbf{y} + \mathbf{v}$. Both algorithms guarantee that \mathbf{w} is close to \mathbf{x} .

Algorithm HERMITE NORMALIZE(C)

Inputs: $C = [c_1 c_2 \cdots c_n]$, the capacities of jugs.

Outputs: $U = [\mathbf{u}_1 \mathbf{u}_2 \cdots \mathbf{u}_n]$ such that $[c_1 c_2 \cdots c_n]U$ is in the Hermite normal form.

Variables: q , temporal storages for $\lfloor \frac{c_i}{c_1} \rfloor$

begin

1. $U := I$;

2. **for** $i = 2$ to n **do**

3. **while** (*true*) **do**

4. $q := \lfloor \frac{c_i}{c_1} \rfloor$; $c_i := c_i - qc_1$; $\mathbf{u}_i := \mathbf{u}_i - q\mathbf{u}_1$;

5. **if** ($c_i = 0$) **then break**;

6. Swap(c_i, c_1); Swap($\mathbf{u}_i, \mathbf{u}_1$);

7. **loop**

8. **next** i

end

Fig. 3. A simple algorithm to compute U .

Theorem 17 ([2]). *The rounding off heuristic algorithm find a vector \mathbf{w} such that $\|\mathbf{x} - \mathbf{w}\|_2 \leq (1 + 2n(\frac{9}{2})^{\frac{1}{2}}) \min_{\mathbf{v} \in L(\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_n)} \|\mathbf{x} - \mathbf{v}\|_2$.*

Theorem 18 ([2]). *The nearest-plane algorithm heuristic algorithm find a vector \mathbf{w} such that $\|\mathbf{x} - \mathbf{w}\|_2 \leq 2^{\frac{n}{2}} \min_{\mathbf{v} \in L(\mathbf{b}_1, \mathbf{b}_2, \dots, \mathbf{b}_n)} \|\mathbf{x} - \mathbf{v}\|_2$.*

Using the nearest-plane heuristic algorithm, we can approximate $\mu_c(x)$ in polynomial time but with an exponential error (approximation ratio).

Corollary 19. *There exists a polynomial-time algorithm to find a vector \mathbf{x} such that $\mathbf{c} \cdot \mathbf{x} = x$ and $\|\mathbf{x}\|_1 \leq \sqrt{n} \cdot 2^{\frac{n-1}{2}} \mu_c(x)$.*

The *LLL* algorithm can reduce a basis of some lattice L to a shorter one. If we replace the *LLL* reduced basis with another shorter basis in the nearest-plane heuristic algorithm, it will have a better performance. Using the basis reduction algorithm in [1], we can have a smaller approximation ratio $2^{O(n \log \log n / \log n)}$.

6. Conclusion and remarks

We have characterized the additively measurable quantities, and proved new lower and upper bounds on the minimum number of measuring steps. We prove that the optimal jug measuring problem is hard to approximate within constant factor. Finally, based on *LLL*-algorithm we give a polynomial-time approximation algorithm with exponential errors.

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