

# Hopf-Galois structures arising from groups with unique subgroup of order $p$ 

Timothy Kohl

For $\Gamma$ a group of order $m p$, where $p$ is a prime with $\operatorname{gcd}(p, m)=1$, we consider the regular subgroups $N \leq \operatorname{Perm}(\Gamma)$ that are normalized by $\lambda(\Gamma)$, the left regular representation of $\Gamma$. These subgroups are in one-to-one correspondence with the Hopf-Galois structures on separable field extensions $L / K$ with $\Gamma=\operatorname{Gal}(L / K)$. Elsewhere we showed that if $p>m$ then all such $N$ lie within the normalizer of the Sylow $p$-subgroup of $\lambda(\Gamma)$. Here we show that one only need assume that all groups of a given order $m p$ have a unique Sylow $p$-subgroup, and that $p$ not be a divisor of the order of the automorphism groups of any group of order $m$. We thus extend the applicability of the program for computing these regular subgroups $N$ and concordantly the corresponding Hopf-Galois structures on separable extensions of degree $m p$.

## Introduction

The motivation and antecedents for this work lie in the subject of Hopf-Galois theory for separable field extensions. Specifically, we extend the results of [Kohl 2013] on Hopf-Galois structures on Galois extensions of degree $m p$ for $p$ a prime where $p>m$. We will not delve into all the particulars of Hopf-Galois theory, since this discussion focuses on the group-theoretic underpinnings of this class of examples. For the general theory, one may consult [Chase and Sweedler 1969] for basic definitions and initial examples and [Greither and Pareigis 1987] for the theory as applied to separable extensions, which is the category in which our earlier paper and others fall. In brief, let $L / K$ be a finite Galois extension with $\Gamma=\operatorname{Gal}(L / K)$. Such an extension is canonically Hopf-Galois for the $K$-Hopf algebra $H=K[\Gamma]$, but also for potentially many other $K$-Hopf algebras. Their enumeration is determined by the following paraphrased variant of the main theorem in [Greither and Pareigis 1987]. (Recall that a regular subgroup of the group of permutations of a set is one whose action on the set is transitive and free.)

[^0]Theorem [Greither and Pareigis 1987]. If $L / K$ is a finite Galois extension with $\Gamma=\operatorname{Gal}(L / K)$ then the Hopf algebras which act to make the extension Hopf-Galois correspond in a one-to-one fashion with regular subgroups $N \leq B=\operatorname{Perm}(\Gamma)$ such that $\lambda(\Gamma) \leq \operatorname{Norm}_{B}(N)$, where $\lambda(\Gamma)$ is the image of the left regular representation of $\Gamma$ in $B$.

Each such $N$ gives rise to the Hopf algebra $H=L[N]^{\Gamma}$, the fixed ring of the group ring $L[N]$ under the action of $\Gamma$ simultaneously on the coefficients, by virtue of $\Gamma=\operatorname{Gal}(L / K)$, and the group elements by virtue of $\lambda(\Gamma)$ normalizing $N$. The problem of enumerating such $N$ for different classes of extensions has been the subject of much recent work by Byott (e.g., [2004]), Childs (e.g., [2003]), the author, and others. This discussion will be strictly focused on the enumeration of these groups, for a particular set of cases, keyed to the order of $\Gamma$ and consequently of any such regular subgroup $N$ which satisfies the conditions of the above theorem. Again, no discussion of Hopf-Galois structures is required from here on since we are looking at the purely group-theoretic translation arising from the above theorem.

## 1. Preliminaries

As it is so essential, let's briefly review regularity in the context of finite groups.
Definition 1.1. Let $X$ be a finite set and let $N \leq \operatorname{Perm}(X)$ be a group of permutations of $X$. We say that $N$ is semiregular if it acts freely, that is, if each element of $N$ apart from the identity acts without fixed points. If $N$ acts transitively on $X$ and $|N|=|X|$ then $N$ is semiregular; and if $N$ is semiregular its action is transitive if and only if $|N|=|X|$. Thus any two of these conditions imply the third. A group satisfying these conditions is called regular.

In view of the cardinality condition, in order to organize the enumeration of the regular $N \leq \operatorname{Perm}(\Gamma)$ that may arise for a given Galois group $\Gamma$, we consider, for [ $M$ ] the isomorphism class of a given group of the same order as $\Gamma$, the set

$$
R(\Gamma,[M])=\left\{N \leq B \mid N \text { regular, } N \cong M, \lambda(\Gamma) \leq \operatorname{Norm}_{B}(N)\right\}
$$

and let $R(\Gamma)$ be the union of the $R(\Gamma,[M])$ over all isomorphism classes of groups of the same order as $\Gamma$.

Since they are important in the enumeration of $R(\Gamma,[M])$, we recall some facts from [Kohl 2013], henceforth referred to as [K].
Lemma 1.2 [K, Corollary 3.3 and Proposition 3.4]. Let $N$ be a regular subgroup of $B$ and define the opposite of $N$ as $N^{\mathrm{opp}}=\operatorname{Cent}_{B}(N)$. Then:

- $N^{\mathrm{opp}}$ is also a regular subgroup of B.
- $N \cap N^{\mathrm{opp}}=Z(N)\left(\right.$ so $N=N^{\mathrm{opp}}$ if and only if $N$ is abelian $)$.
- $\left(N^{\mathrm{opp}}\right)^{\mathrm{opp}}=N$.
- $N \in R(\Gamma,[M])$ if and only if $N^{\mathrm{opp}} \in R(\Gamma,[M])$.

Again, in the cases considered in $[\mathrm{K}]$ it was assumed that $|\Gamma|=m p$ for $p$ prime and $p>m$. Our goal is to extend those results to groups of order $m p$ where $\operatorname{gcd}(p, m)=1$, but where one need not assume that $p>m$.

We begin by briefly reviewing the setup in [K], where we considered groups $\Gamma$ of order $m p$ with a unique and therefore characteristic Sylow $p$-subgroup due to the assumption that $p>m$. Since $p>m$ obviously implies $\operatorname{gcd}(p, m)=1$, by the Schur-Zassenhaus lemma $\Gamma$ may be written as $P Q$ for $P$ and $Q$ subgroups of $\Gamma$ where $|P|=p$ and $|Q|=m$. More specifically, there is a split exact sequence $P \rightarrow \Gamma \rightarrow Q$ whereby $\Gamma=P \rtimes_{\tau} Q$, with $\tau: Q \rightarrow \operatorname{Aut}(P)$ induced by conjugation within $\Gamma$ by the complementary subgroup $Q$. Using $Q$ for the quotient of $\Gamma$ by $P$ and the image of the section in $\Gamma$ is admittedly a slight abuse of notation. The condition $p>m$ is sufficient, of course, to make the Sylow $p$-subgroup unique and have order $p$.

Going forward, we wish to drop the assumption that $p>m$ and consider groups of order $m p$ for $p$ prime, with $\operatorname{gcd}(p, m)=1$ and where congruence conditions force any group of order $m p$, including $\Gamma$ and any $N \in R(\Gamma)$, to have a unique Sylow p-subgroup.

With $\Gamma$ as above, if $\lambda: \Gamma \rightarrow \operatorname{Perm}(\Gamma)=B$ is the left regular representation then we define $\mathcal{P}=P(\lambda(\Gamma))$ to be the Sylow $p$-subgroup of $\lambda(\Gamma)$ and $\mathcal{Q}$ to be the complementary subgroup to $\mathcal{P}$ in $\lambda(\Gamma)$. We wish to prove the following strengthening of [K, Theorem 3.5] which will allow us to apply the program developed in Sections $1-3$ of $[\mathrm{K}]$ (and applied in subsequent sections therein) to a much larger class of groups.

Theorem 1.3. Let $\Gamma$ have order $m p$ where $\operatorname{gcd}(p, m)=1$ and $p \nmid|\operatorname{Aut}(Q)|$ for any group $Q$ of order $m$, and assume any group of order mp has a unique Sylow p-subgroup. If $N \in R(\Gamma)$ then $N$ is a subgroup of $\operatorname{Norm}_{B}(\mathcal{P})$.

To prove this, we need to modify certain key results from [K], starting with Lemma 1.1 regarding the $p$-torsion of $\operatorname{Aut}(\Gamma)$.

Lemma 1.4. Let $\Gamma$ have order $m p$, where $\operatorname{gcd}(p, m)=1$, and assume $\Gamma$ has a unique Sylow p-subgroup, so that $\Gamma \cong P \rtimes_{\tau} Q$ as above. Assume also that $p$ does not divide $|\operatorname{Aut}(Q)|$.
(a) If $\tau$ is trivial (whence $\Gamma \cong P \times Q$ ) then $p$ does not divide $|\operatorname{Aut}(\Gamma)|$.
(b) If $\tau$ is nontrivial then $\operatorname{Aut}(\Gamma)$ has a unique Sylow p-subgroup, consisting of the inner automorphisms induced by conjugation by elements of $P$.

Proof. In (a), if $\Gamma$ is such a direct product then $\operatorname{Aut}(\Gamma)=\operatorname{Aut}(P) \times \operatorname{Aut}(Q)$ and so if $p \nmid|\operatorname{Aut}(Q)|$ then $p \nmid|\operatorname{Aut}(\Gamma)|$. The proof is basically the same as in [K]. For (b),
since $\Gamma / P \cong Q$, any $\psi \in \operatorname{Aut}(\Gamma)$ induces $\bar{\psi} \in \operatorname{Aut}(\Gamma / P) \cong \operatorname{Aut}(Q)$ and if $\psi$ has $p$-power order then $\bar{\psi}$ is trivial since $p$ does not divide $|\operatorname{Aut}(Q)|$. And, as also observed in $[\mathrm{K}]$, when $\Gamma$ is not a direct product, $|P \cap Z(\Gamma)|=1$ and conjugation in $\Gamma$ by elements of $P$ yields the order- $p$ subgroup of $\operatorname{Aut}(\Gamma)$.

The condition that $p \nmid|\operatorname{Aut}(Q)|$ was automatic when $p>m$, but it often holds true even when $p<m$. For example, if $p=5$ and $m=8$ then Sylow theory easily shows that any group of order 40 will have a unique Sylow $p$-subgroup. One could also consider the groups of order eight $-C_{8}, C_{4} \times C_{2}, C_{2} \times C_{2} \times C_{2}, D_{4}$ and $Q_{2}$ whose automorphism groups have orders $4,8,168,8,24$, respectively, none of which is divisible by 5 .

The cycle structure of a regular permutation group's elements is greatly circumscribed by the condition that all nontrivial elements of the group act freely. Any such element, because it and all its nontrivial powers lack fixed points, must be a product of cycles of the same length, and the sum of the lengths must equal $|X|$. For example, if $X=\{1,2,3,4,5,6\}$ then $(1,2)(3,4)(5,6)$ and $(1,2,3)(4,5,6)$ satisfy this property. In contrast, $\mu=(1,2,3,4)(5,6)$ cannot belong to a regular subgroup of $\operatorname{Perm}(X)$ even though it does not have fixed points, because $\mu^{2}=(1,3)(2,4)$ does.

Since $\mathcal{P}$ is a nontrivial subgroup of the canonically regular permutation group $\lambda(\Gamma)$, we must have

$$
\mathcal{P}=\langle\pi\rangle=\left\langle\pi_{1} \pi_{2} \cdots \pi_{m}\right\rangle,
$$

where the $\pi_{i}$ are disjoint $p$-cycles. In a similar fashion, if $N$ is any regular subgroup of $B$ then its Sylow $p$-subgroup $P(N)$ is also cyclic of order $p$ and therefore of the form

$$
P(N)=\langle\theta\rangle=\left\langle\theta_{1} \theta_{2} \cdots \theta_{m}\right\rangle,
$$

where the $\theta_{i}$ are also disjoint $p$-cycles. For $N \in R(\Gamma)$ we are looking at those regular $N$ which are normalized by $\lambda(\Gamma)$. Now, $P(N)$ is characteristic in $N$, so $\lambda(\Gamma)$ normalizing $N$ implies that $\lambda(\Gamma)$ (and therefore $\mathcal{P}$ ) normalizes $P(N)$. In $[\mathrm{K}]$ the assumption $p>m$ was used to show that $\mathcal{P}$ and $P(N)$ must, in fact, centralize each other. In particular Proposition 1.2 there showed that if $p>m$ then (after renumbering the $\theta_{i}$ if necessary) one has $\theta_{i}=\pi_{i}^{a_{i}}$ for $a_{i} \in \mathbb{F}_{p}^{\times}$. The reason for this was that for $p>m$ the group $S_{m}$ contains no elements of order $p$, and so $\theta$ is a product of the same $\pi_{i}$ that comprise the generator of $\mathcal{P}$.

As it turns out, this is not automatically true if we just assume that $\operatorname{gcd}(p, m)=1$. For example, if $p=5$ and $m=8$ then in $S_{40}$ let

$$
\pi_{i}=(1+(i-1) 5,2+(i-1) 5,3+(i-1) 5,4+(i-1) 5,5+(i-1) 5)
$$

for $i=1, \ldots, 8$ and let $\theta_{j}=(j, j+5, j+10, j+15, j+20)$ for $j=1, \ldots, 5$ and $\theta_{6}=\pi_{6}, \theta_{7}=\pi_{7}, \theta_{8}=\pi_{8}$. One may verify that $\pi=\pi_{1} \cdots \pi_{8}$ is centralized by $\theta=\theta_{1} \cdots \theta_{8}$ but $\theta_{j}$, for $j=1, \ldots, 5$, is not a power of any $\pi_{i}$.

This example shows that the $P(N) \leq N$ being normalized, and thus centralized, by $\mathcal{P}$ is insufficient to guarantee that $P(N) \leq\left\langle\pi_{1}, \pi_{2}, \ldots, \pi_{m}\right\rangle$. This does not nullify the possibility of the program in $[\mathrm{K}]$ being generalized. This example merely shows that $\operatorname{Cent}_{B}(\mathcal{P})$ contains many semiregular subgroups of order $p$ that are not subgroups of $\left\langle\pi_{1}, \pi_{2}, \ldots, \pi_{m}\right\rangle$. However, it turns out that for those $N$ normalized by $\lambda(\Gamma)$, since $P(N)$ is characteristic in $N$ and therefore normalized by $\lambda(\Gamma)$, the possible $P(N)$ that can arise are still contained in $\left\langle\pi_{1}, \pi_{2}, \ldots, \pi_{m}\right\rangle$. To arrive at this conclusion, we need to recall some facts about the structure of $\operatorname{Norm}_{B}(\mathcal{P})$ and $\operatorname{Cent}_{B}(\mathcal{P})$.

With $\mathcal{P}=\langle\pi\rangle=\left\langle\pi_{1} \cdots \pi_{m}\right\rangle$, where the $\pi_{i}$ are disjoint $p$-cycles, we can define $\mathcal{V}=\left\langle\pi_{1}, \pi_{2}, \ldots, \pi_{m}\right\rangle$, the elementary abelian subgroup of $B$ generated by the $\pi_{i}$. Also, we can choose $\gamma_{i} \in \Gamma$ for $i=1, \ldots, m$ such that $\pi_{i}=\left(\gamma_{i}, \pi\left(\gamma_{i}\right), \ldots, \pi^{p-1}\left(\gamma_{i}\right)\right)$ and if we let $\Pi_{i}$ be the support of $\pi_{i}$, then the $\Pi_{i}$ are, of course, disjoint and their union is $\Gamma$ as a set. Define $\mathcal{S} \leq B$ to be those permutations $\bar{\alpha}$ such that for each $i \in\{1, \ldots, m\}$ there exists a single $j \in\{1, \ldots, m\}$ satisfying $\bar{\alpha}\left(\pi^{t}\left(\gamma_{i}\right)\right)=\pi^{t}\left(\gamma_{j}\right)$ for each $t \in \mathbb{Z}_{p}$. Equivalently, $\bar{\alpha}$ operates on the blocks $\Pi_{i}$ as follows:

$$
\bar{\alpha}\left(\left\{\gamma_{i}, \pi\left(\gamma_{i}\right), \ldots, \pi^{p-1}\left(\gamma_{i}\right)\right\}\right)=\left\{\gamma_{j}, \pi\left(\gamma_{j}\right), \ldots, \pi^{p-1}\left(\gamma_{j}\right)\right\}
$$

It is clear that $\mathcal{S}$ is isomorphic to $S_{m}$ viewed as $\operatorname{Perm}\left(\left\{\Pi_{1}, \ldots, \Pi_{m}\right\}\right)$, where $\bar{\alpha} \in \mathcal{S}$ corresponds to a permutation $\alpha \in S_{m}$ which permutes the $m$ blocks $\Pi_{i}$ amongst each other. In a similar fashion, we may define another subgroup $\mathcal{U} \leq B$ keyed to $\pi$ and the $\pi_{i}$. For a cyclic group $C=\langle x\rangle$ of order $p$, the automorphisms are given by $x \mapsto x^{c}$ for $c \in U_{p}=\mathbb{F}_{p}^{\times}=\langle u\rangle$. Within $B$, therefore, since $\mathcal{P}$ is cyclic of order $p$, there exists a product $u_{1} \cdots u_{m}$ of $m$ disjoint ( $p-1$ )-cycles with the property that $u_{i} \pi_{i} u_{i}^{-1}=\pi_{i}^{u}$, and $u_{i} \pi_{j} u_{i}^{-1}=\pi_{j}$ for $j \neq i$. (Note that the support of each $u_{i}$ is $\Pi_{i}-\{$ one point $\}$.) Therefore, $\left(u_{1} \cdots u_{m}\right) \pi\left(u_{1} \cdots u_{m}\right)^{-1}=\pi^{u}$ and we define $\mathcal{U}=\left\langle u_{1} \cdots u_{m}\right\rangle$. With these three subgroups of $B$ so defined, we can easily describe $\operatorname{Cent}_{B}(\mathcal{P})$ and $\operatorname{Norm}_{B}(\mathcal{P})$ as in [K] by

$$
\begin{aligned}
\operatorname{Cent}_{B}(\mathcal{P}) & =\mathcal{V} \mathcal{S} \cong C_{p} 2 S_{m} \cong C_{p}^{m} \rtimes S_{m} \\
\operatorname{Norm}_{B}(\mathcal{P}) & =\mathcal{V} \mathcal{U} \cong C_{p}^{m} \rtimes\left(\operatorname{Aut}\left(C_{p}\right) \times S_{m}\right)
\end{aligned}
$$

where $C_{p}$ denotes the cyclic group of order $p$ and $S_{m}$ is the $m$-th symmetric group.
The semidirect product formulation is useful and may be closely connected to the intrinsic (as a subgroup of $B$ ) description. We may view $\mathcal{V}=\left\langle\pi_{1}, \ldots, \pi_{m}\right\rangle$ naturally as $C_{p}^{m}$ but also, more fruitfully, as $\mathbb{F}_{p}^{m}$, the dimension- $m$ vector space over $\mathbb{F}_{p}$, so that we may equate $\pi_{1}^{a_{1}} \cdots \pi_{m}^{a_{m}}$ with $\left[a_{1}, \ldots, a_{m}\right]$, a vector in $\mathbb{F}_{p}^{m}$. As the group $\mathcal{S}$ permutes the $\pi_{i}$ amongst themselves, we may identify it with permutations $\alpha \in S_{m}$ acting by coordinate shifts on the vectors $\hat{a}=\left[a_{1}, \ldots, a_{m}\right]$ and where $u \in U_{p}$ acts by scalar multiplication. Thus we may represent a typical element of $\operatorname{Norm}_{B}(\mathcal{P})$ by a triple $\left(\hat{a}, u^{r}, \alpha\right)$, with $\hat{a} \in \mathbb{F}_{p}^{m}, u \in U_{p}$ and $\alpha \in S_{m}$, where (as permutations)
$\left(\hat{a}, u^{r}, \alpha\right)\left(\pi_{i}^{k}\left(\gamma_{i}\right)\right)=\pi_{\alpha(i)}^{k u^{r}+a_{\alpha(i)}}\left(\gamma_{\alpha(i)}\right)$ and where the multiplication (and resulting conjugation operations) are defined by

$$
\begin{align*}
\left(\hat{a}, u^{r}, \alpha\right)\left(\hat{b}, u^{s}, \beta\right) & =\left(\hat{a}+u^{r} \alpha(\hat{b}), u^{r+s}, \alpha \beta\right),  \tag{1}\\
\left(\hat{b}, u^{s}, \beta\right)\left(\hat{a}, u^{r}, \alpha\right)\left(\hat{b}, u^{s}, \beta\right)^{-1} & =\left(\hat{b}+u^{s} \beta(\hat{a})-u^{r}\left(\beta \alpha \beta^{-1}\right)(\hat{b}), u^{r}, \beta \alpha \beta^{-1}\right),  \tag{2}\\
\left(\hat{a}, u^{r}, \alpha\right)^{n} & =\left(\sum_{t=0}^{n-1} u^{r t} \alpha^{t}(\hat{a}), u^{r n}, \alpha^{n}\right) . \tag{3}
\end{align*}
$$

In this setting the elements of $\mathcal{V}$ correspond to tuples of the form $(\hat{v}, 1, I)$, where $I$ is the identity of $S_{m}$; in particular $\pi=\pi_{1} \pi_{2} \cdots \pi_{m}=([1,1, \ldots, 1], 1, I)$. The elements of $\operatorname{Cent}_{B}(\mathcal{P})$ correspond to those tuples where $r=0$ (i.e., the middle coordinate is 1 ), which leads us back to the discussion of $P(N)$ for $N$ a regular subgroup of $B$ normalized by $\lambda(\Gamma)$. In this situation we have $P(N)=\langle\theta\rangle$, where $\theta=(\hat{a}, 1, \alpha)$ has order $p$ and no fixed points. If $P(N) \notin \mathcal{V}$ then $\alpha \neq I$, implying (since $\alpha \in S_{m}$ with $m$ coprime to $p$ ) that $\alpha$ has fixed points in $\{1, \ldots, m\}$. If $\alpha(i)=i$ then

$$
\theta\left(\pi_{i}^{k}\left(\gamma_{i}\right)\right)=\pi_{\alpha(i)}^{k+a_{\alpha(i)}}\left(\gamma_{\alpha(i)}\right)=\pi_{i}^{k+a_{i}}\left(\gamma_{i}\right),
$$

which means that $a_{i} \neq 0$ and more importantly that $\theta$ restricted to $\Pi_{i}$ equals $\pi_{i}^{a_{i}}$. And for those $j$ not fixed by $\alpha$, the restriction of $\theta$ to $\Pi_{j}$ is not a power of $\pi_{j}$. That is, $\theta=\theta_{1} \theta_{2} \cdots \theta_{m}$, where $\theta_{i}=\pi_{i}^{a_{i}}$ for at least one $i$, and $\theta_{j} \notin \mathcal{V}$ for at least one $j$. The example given above for $S_{40}$ is an instance of this; in particular, the fixed-point-free element of order 5 is ( $[1,1,1,1,1,1,1,1], 1,(1,2,3,4,5)$ ), which is in Cent $S_{40}\left(\left\langle\pi_{1} \pi_{2} \cdots \pi_{8}\right)\right\rangle$.

The above example motivates the following.
Definition 1.5. For $\theta \in B$ and $\pi_{i}$ as above, we say $\pi_{i}$ divides $\theta$ (denoted $\left.\pi_{i} \mid \theta\right)$ if the cycle structure of $\theta$ contains some nontrivial power of $\pi_{i}$. Similarly we write $\pi_{i} \nmid \theta$ if no power of $\pi_{i}$ is a factor in the cycle structure of $\theta$.

Observe that $\pi_{i} \mid \theta$ if and only if $\pi_{i} \mid \theta^{e}$ for all $e \in U_{p}$.
The requirement that $N$ be normalized by $\lambda(\Gamma)$, together with the fact that $P(N)$ is characteristic, means that $P(N)$ is normalized by $\lambda(\Gamma)$. The upshot of this is the following recapitulation of [K, Proposition 1.2], which (together with Lemma 1.4) will allow us to deduce our main result, namely that $N \in R(\Gamma)$ implies $N \leq \operatorname{Norm}_{B}(\mathcal{P})$ under weaker hypotheses:
$p$ and $m$ are coprime, $p$ does not divide $|\operatorname{Aut}(Q)|$ for any group $Q$ of order $m$, and any group of order $m p$ has a unique Sylow $p$-subgroup.

This is the core result, since it guarantees that $P(N) \leq \mathcal{V} \leq \operatorname{Norm}_{B}(\mathcal{P})$ for all $N \in R(\Gamma)$.

Theorem 1.6. If $N$ is a regular subgroup of $B$ normalized by $\lambda(\Gamma)$ and $P(N)$ is its Sylow $p$-subgroup, then $P(N)$ is a semiregular subgroup of $\mathcal{V}=\left\langle\pi_{1}, \pi_{2}, \ldots, \pi_{m}\right\rangle$.
That is, $P(N)=\left\langle\pi_{1}^{a_{1}} \cdots \pi_{m}^{a_{m}}\right\rangle$, where $a_{i} \in U_{p}=\mathbb{F}_{p}^{\times}$for $i=1, \ldots, m$.
Proof. If $P(N)=\langle\theta\rangle$ is not a subgroup of $\mathcal{V}$ then, as shown above, $\theta=\theta_{1} \theta_{2} \cdots \theta_{m}$, where $\pi_{i} \mid \theta$ for some $i$ and $\pi_{j} \nmid \theta$ for some $j \neq i$. In [K, Proposition 3.8] it is shown that if $\left(\hat{b}, u^{s}, \beta\right) \in \operatorname{Norm}_{B}(\mathcal{P})$ has order coprime to $p$ then the permutation coordinate $\beta$ acts without fixed points. It is important to note that the proof of this is not dependent on whether $p<m$ or $p>m$. Applying this to $\lambda(\Gamma)$, which normalizes $\mathcal{P}=\left\langle\pi_{1} \cdots \pi_{m}\right\rangle$, has no fixed points, and contains elements $\left(\hat{b}, u^{s}, \beta\right)$ of order coprime to $p$, we conclude that $\beta \in S_{m}$ is fixed-point-free. In fact, if $\mathcal{Q}$ is the complementary subgroup of $\mathcal{P}$ in $\lambda(\Gamma)$ and if $t$ maps $\left(\hat{b}, u^{s}, \beta\right)$ to $\beta$ then $t(\mathcal{Q})$ must be a regular subgroup of $S_{m}$. The reason is that the elements of $\mathcal{Q}$ have order relatively prime to $p$, so $t(\mathcal{Q})$ is a semiregular subgroup of $S_{m}$; but if $|t(\mathcal{Q})|<m$ there would exist $\left(\hat{a}, u^{r}, \alpha\right)$ and $\left(\hat{a}^{\prime}, u^{r^{\prime}}, \alpha\right)$ in $\mathcal{Q}$, which would imply that $\left(\hat{a}^{\prime}-u^{r^{\prime}-r} \hat{a}, u^{r^{\prime}-r}, I\right) \in \mathcal{Q}$, which, again by $[\mathrm{K}$, Proposition 3.8], must have order divisible by $p$. But since this element is in $\mathcal{Q}$ it must be the identity, whence $r=r^{\prime}$ and so $\hat{a}=\hat{a}^{\prime}$. This means that $t(\mathcal{Q})$ is regular; by transitivity we pick an element $g=\left(\hat{b}, u^{s}, \beta\right)$ in $\mathcal{Q}$ with $\beta(i)=j$, and then $g([1,1, \ldots, 1], 1, I) g^{-1}=$ $\left(u^{s} \beta([1,1, \ldots, 1]), 1, I\right)=\left(u^{s}[1,1, \ldots, 1], 1, I\right)$, where, in particular,

$$
g \pi_{1} \pi_{2} \cdots \pi_{m} g^{-1}=\pi_{\beta(1)}^{u^{s}} \pi_{\beta(2)}^{u^{s}} \cdots \pi_{\beta(m)}^{u^{s}} .
$$

And since $g\left(\theta_{1} \theta_{2} \cdots \theta_{m}\right) g^{-1}=\left(g \theta_{1} g^{-1}\right)\left(g \theta_{2} g^{-1}\right) \cdots\left(g \theta_{m} g^{-1}\right)$, we have $g \theta_{i} g^{-1}=$ $g \pi_{i}^{a_{i}} g^{-1}=\pi_{\beta(i)}^{u^{s} a_{i}}=\pi_{j}^{u^{s} a_{i}}$; therefore $\pi_{j} \mid g \theta g^{-1}$. The problem now is that $g \theta g^{-1}=\theta^{e}$ for some $e \in U_{p}$ implies that $\pi_{j} \mid \theta$, contrary to the assumption that $\pi_{j} \nmid \theta$. We conclude that any such $\theta$ must, in fact, be a fixed-point-free subgroup of $\mathcal{V}$ and therefore of the form asserted in the statement of the theorem.

With this in place, we can review the remaining foundational elements in [K], which, without serious modifications, imply, under the weaker hypotheses (4), that $N \leq \operatorname{Norm}_{B}(\mathcal{P})$ for all $N \in R(\Gamma)$. We do this also in order to provide some applications for classes of groups where these weaker hypotheses hold.

We have just shown that $N \in R(\Gamma)$ implies $P(N) \leq \operatorname{Norm}_{B}(\mathcal{P})$ and that $P(N)=$ $\left\langle\pi_{1}^{a_{1}} \cdots \pi_{m}^{a_{m}}\right\rangle$. Define $Q(N)$ to be the complementary subgroup to $P(N)$ inside $N$. Then, since $Q(N)$ normalizes $P(N)$, we have $q \pi_{i}^{a_{i}} q^{-1}=\pi_{j}^{b_{j}}$ for $q \in Q(N)$, where the mapping of $i \mapsto j=q(i)$ for $i, j \in\{1, \ldots, m\}$ makes $Q(N)$ (abstractly) a regular subgroup of $S_{m}$. Let $\mathcal{Q}=Q(\lambda(\Gamma))$ be the complementary subgroup to $\mathcal{P}$ inside $\lambda(\Gamma)$. If $N \in R(\Gamma)$, then the elements of $\lambda(\Gamma)$ that act nontrivially on $P(N)$ are those in $\mathcal{Q}$. If we define $\hat{v}_{i}=[0, \ldots, 1, \ldots, 0]=\pi_{i}$ then we have the following result (whose proof does not require that $p<m$ ) about the possibilities for $P(N)$ for any $N \in R(\Gamma)$.

Theorem 1.7 [K, Theorem 2.1]. Any semiregular subgroup of $\mathcal{V}$ of order $p$ that is normalized by $\mathcal{Q}$ is generated by

$$
\hat{p}_{\chi}=\sum_{\gamma \in \mathcal{Q}} \chi(\gamma) \hat{v}_{\gamma(1)}
$$

where $\chi: \mathcal{Q} \rightarrow U_{p}=\mathbb{F}_{p}^{*}$ is a linear character of $\mathcal{Q}$.
This theorem allows one to compute "potential" $P(N)$, one for each such linear character. For example, $\mathcal{P}$ is generated by $[1,1, \ldots, 1]=\hat{p}_{\iota}$ where $\iota: \mathcal{Q} \rightarrow \mathbb{F}_{p}^{*}$ is the trivial character. All of these order- $p$ elements have the form $\left(\hat{p}_{\chi}, 1, I\right)$ in the semidirect product formulation of $\operatorname{Norm}_{B}(\mathcal{P})$. The remaining component is to show that not only does $N \in R(\Gamma)$ imply $P(N) \leq \operatorname{Norm}_{B}(\mathcal{P})$ but also that $N$ itself is contained in $\operatorname{Norm}_{B}(\mathcal{P})$. The assumption $p>m$ in [K] was not actually used in the proof of this main result, but rather in the following lemma:

Lemma 1.8 [K, Lemma 2.2]. Let $\chi_{1}, \chi_{2}$ be distinct linear characters of $\mathcal{Q}$ in $\mathbb{F}_{p}^{*}$. Then $\left\langle\hat{p}_{\chi_{1}}, \hat{p}_{\chi_{2}}\right\rangle$ cannot contain $\hat{p}_{l}$.

A careful reading of the proof of this in [K] shows it is not necessary to assume that $p>m$, but merely $\operatorname{gcd}(p, m)=1$. With this completed, Theorem 1.3 , saying that $N \in R(\Gamma)$ implies $N \leq \operatorname{Norm}_{B}(\mathcal{P})$, follows in the exact same fashion as in [K], since the proof does not hinge on the relationship between $p$ and $m$ beyond the fact that they are relatively prime. It does require that $p$ not divide the order of $\operatorname{Aut}(Q)$, as in Lemma 1.4. This assumption on $|\operatorname{Aut}(Q)|$ is needed to control the size and structure of the Sylow $p$-subgroup of $\operatorname{Norm}_{B}(N)$. Specifically, it is either cyclic of order $p$ if $P(N)$ is central in $N$, or elementary abelian of order $p^{2}$ if $P(N)$ is noncentral. Again, in [K] this was automatic from assuming that $p>m$.

The application of this theorem, which is the actual program in [K] (demonstrated in Sections 4 and 5 therein), is based on the observation (in [K, Proposition 3.11]) that any two regular subgroups of $B$ that are isomorphic as abstract groups are, in fact, conjugate subgroups of $B$. That being said, to enumerate $R(\Gamma)$, one can avoid the complications of working with left regular representations and instead:
(1) replace $B=\operatorname{Perm}(\Gamma)$ by $S_{m p}=\operatorname{Perm}(\{1, \ldots, m p\})$,
(2) choose $\mathcal{P}=\left\langle\pi_{1} \cdots \pi_{m}\right\rangle$ where $\pi_{i}=(1+(i-1) p, \ldots, p i)$,
(3) determine $\mathcal{Q}$ corresponding to each such $\Gamma$ where $\Gamma=\mathcal{P} \mathcal{Q}$,
(4) embed the $\Gamma$ as subgroups of the semidirect product formulation of $\operatorname{Norm}_{S_{m p}}(\mathcal{P})$,
(5) enumerate the characters $\chi: \mathcal{Q} \rightarrow \mathbb{F}_{p}^{*}$ and concordantly the potential $P(N)$ as $\left\langle\hat{p}_{\chi}\right\rangle$ also embedded in $\operatorname{Norm}_{S_{m p}}(\mathcal{P})$,
(6) compute the possible $N \in R(\Gamma)$ that may arise, also as subgroups of $\operatorname{Norm}_{S_{m p}}(\mathcal{P})$.

If one is more interested in the sizes of the different $R(\Gamma,[M])$ one may use the fact (in [K, Theorem 3.5]) that, for each $N \in R(\Gamma,[M])$, if $P(N)$ is central in $N$ then $P(N)=P\left(N^{\mathrm{opp}}\right)$, and otherwise either $P(N)=\mathcal{P}$ or $P\left(N^{\mathrm{opp}}\right)=\mathcal{P}$. The point is that one can enumerate those $N \in R(\Gamma,[M])$ for which $P(N)=\mathcal{P}$ and, depending on whether $P(N)$ is central, use the above fact to infer the count of those for $N$ for which $P(N) \neq \mathcal{P}$ (if any). The virtue of this is that one need not calculate the characters of $\mathcal{Q}$ in $\mathbb{F}_{p}^{*}$, nor the resulting potential $P(N)$.

We shall demonstrate applications of this program, where we now have a wider class of examples to choose from, based upon the conditions on $p, m$ and $|\operatorname{Aut}(Q)|$ as discussed above.

## 2. Groups of order $\boldsymbol{p}_{1} \boldsymbol{p}_{2} \boldsymbol{p}_{3}$

To be slightly formal, if $n_{p}$ denotes the number of Sylow $p$-subgroups of a group, we define the following subsets of $\mathbb{N} \times \mathbb{N}$ :
$F_{Q}=\{(p, m) \mid p$ prime, $\operatorname{gcd}(p, m)=1, p \nmid|\operatorname{Aut}(Q)|$ for all groups $Q$ of order $m\}$, $F_{S}=\left\{(p, m) \mid p\right.$ prime, $\operatorname{gcd}(p, m)=1, n_{p}=1$ for all groups of order $\left.m p\right\}$.

The program in $[\mathrm{K}]$ for enumerating Hopf-Galois structures on Galois extensions of order $m p$ may be used for those $(p, m) \in F_{Q} \cap F_{S}$. As in [K], $(p, m) \in F_{Q} \cap F_{S}$ for $p$ prime when $p>m$, but we want to now consider other $p$ and $m$. The case of $p=5$ and $m=8$ as indicated already is one such example. In lieu of working out the enumeration of all the $14^{2}$ possible pairings $R(\Gamma,[M])$ for order 40 , we shall instead conclude with an overview of some (classes of) choices for $|\Gamma|=|N|=n=p m$ which force $(p, m) \in F_{Q} \cap F_{S}$. Such forcing conditions have appeared in group theory literature including recent examples such as [Pakianathan and Shankar 2000]. Our example will be somewhat more narrow, but is in this same spirit.

If $p_{1}<p_{2}<p_{3}$ are primes, then by Sylow theory $n_{p_{3}} \equiv 1\left(\bmod p_{3}\right)$ and $n_{p_{3}} \mid p_{1} p_{2}$. However, $n_{p_{3}} \neq p_{1}$ and $n_{p_{3}} \neq p_{2}$ since $p_{3}$ is larger than $p_{1}$ and $p_{2}$. If $n_{p_{3}}=p_{1} p_{2}$ then one must have $p_{1} p_{2}\left(p_{3}-1\right)$ elements of order $p_{3}$ and so, by necessity, $n_{p_{2}}=1$. Thus $\Gamma$ has a normal subgroup $P_{2}$ of order $p_{2}$ and so $\Gamma / P_{2}$ (having order $p_{1} p_{3}$ ) has a normal subgroup of order $p_{3}$, which gives rise to a normal abelian subgroup $\Delta \leq \Gamma$ of order $p_{2} p_{3}$. We have that $P_{2} \leq \Delta$ but also that $\Delta$ must contain a normal (in particular characteristic) subgroup $P_{3}$ of order $p_{3}$, which means that $P_{3} \triangleleft \Gamma$, so that, in fact, $n_{p_{3}}=1$. Hence ( $p_{3}, p_{1} p_{2}$ ) is guaranteed to be in $F_{S}$. Moreover, the complementary subgroup $Q$ is either a cyclic or metacyclic group of order $p_{1} p_{2}$. The question then is whether $\left(p_{3}, p_{1} p_{2}\right) \in F_{Q}$ as well. However, this is easy since

$$
|\operatorname{Aut}(Q)|= \begin{cases}\left(p_{1}-1\right)\left(p_{2}-1\right) & \text { if } Q \text { is abelian, } \\ p_{2}\left(p_{2}-1\right) & \text { if } Q \text { is nonabelian },\end{cases}
$$

and so, if $p_{1}<p_{2}<p_{3}$ one has $\left(p_{3}, p_{1} p_{2}\right) \in F_{S} \cap F_{Q}$.

If $|\Gamma|=p_{1} p_{2} p_{3}$ then $\Gamma=P Q$ where $P$ is cyclic of order $p_{3}$, of course, and normalized by $Q$ which has order $p_{1} p_{2}$, which is itself either a direct or semidirect product of cyclic groups. That is, $\Gamma \cong C_{p_{3}} \rtimes_{g}\left(C_{p_{2}} \rtimes_{f} C_{p_{1}}\right)$ where $f: C_{p_{1}} \rightarrow \operatorname{Aut}\left(C_{p_{2}}\right)$ and $g: C_{p_{2}} \rtimes_{f} C_{p_{1}} \rightarrow \operatorname{Aut}\left(C_{p_{3}}\right)$. There are at most two groups $Q$ of order $p_{1} p_{2}$ depending on whether $f$ is nontrivial. The group $\Gamma$ being an iterated (semi)direct product, the role of $g$ must also be factored into the enumeration of the distinct groups of order $p_{1} p_{2} p_{3}$. In particular, we must consider whether the $C_{p_{2}}$ and/or the $C_{p_{1}}$ components of $Q$ act nontrivially on $C_{p_{3}}$. These possibilities for $f$ and $g$ are keyed to congruence conditions on the $p_{i}$, in particular, whether $p_{1} \mid\left(p_{2}-1\right)$ and/or $p_{1} \mid\left(p_{3}-1\right)$ and/or $p_{2} \mid\left(p_{3}-1\right)$. Alonso [1976] (while exploring an explicit formula due to Hölder [1895] for the number of groups of square-free order) works through the enumeration of groups having order equal to the product of three distinct primes. In particular we give Table 1 therein of the number of groups of order $p_{1} p_{2} p_{3}$ (with our notation for the three primes):

| $p_{2} \mid\left(p_{3}-1\right)$ | $p_{1} \mid\left(p_{3}-1\right)$ | $p_{1} \mid\left(p_{2}-1\right)$ | \# groups |
| :---: | :---: | :---: | :---: |
| no | no | no | 1 |
| no | no | yes | 2 |
| no | yes | no | 2 |
| no | yes | yes | $p_{1}+2$ |
| yes | no | no | 2 |
| yes | no | yes | 3 |
| yes | yes | no | 4 |
| yes | yes | yes | $p_{1}+4$ |

For two of the eight cases, the number of groups varies linearly with $p_{1}$ (specifically when $C_{p_{1}}$ acts nontrivially on both $C_{p_{2}}$ and $C_{p_{3}}$ ) but for the others the size is constant. For the case of $p_{2}\left|\left(p_{3}-1\right), p_{1}\right|\left(p_{3}-1\right)$, and $p_{1} \mid\left(p_{2}-1\right)$, it follows that $p_{3}>p_{1} p_{2}=m$ which falls into the category of cases dealt with in [K]. Indeed, therein we enumerated $R(\Gamma,[M])$ for all groups of order $m p$ where $p=2 q+1$ for $q$ a prime (making $p$ a safe-prime) and $m=\phi(p)=2 q$ so that $m p=2 \times q \times(2 q+1)$, the product of three distinct primes! Therefore we will instead consider the case where $p_{3} \not \equiv 1\left(\bmod p_{2}\right)$ but where $p_{1}$ divides both $p_{2}-1$ and $p_{3}-1$, for this includes cases where $p_{3}<m=p_{1} p_{2}$, for example, $\left(p_{1}, p_{2}, p_{3}\right)=(3,7,13)$. The $p_{1}+2$ cases can be presented explicitly and again we refer to [Alonso 1976] for the particulars with a slight modification of his notation.

Given $\left(p_{1}, p_{2}, p_{3}\right)$, the groups of order $p_{1} p_{2} p_{3}$ are iterated semidirect products, the number of which, as mentioned above, are keyed to elements of order $p_{1}$ in $U_{p_{2}}$ and $U_{p_{3}}$. Specifically, if $U_{p_{3}}=\left\langle u_{3}\right\rangle$ and $U_{p_{2}}=\left\langle u_{2}\right\rangle$, then the conditions

$$
\left.p_{3} \equiv 1\left(\bmod p_{1}\right), \quad p_{2} \equiv 1\left(\bmod p_{1}\right)\right), \quad v_{3}=u_{3}^{\left(p_{3}-1\right) / p_{1}}, \quad v_{2}=u_{2}^{\left(p_{2}-1\right) / p_{1}}
$$

together imply that then $\left|v_{3}\right|=p_{1}$ and $\left|v_{2}\right|=p_{1}$. We have the following presentations for groups of order $p_{1} p_{2} p_{3}$ which are the cases listed as (10)-(13) in [Alonso 1976, p. 634]. (Note that Alonso adopts the ordering $p_{3}<p_{2}<p_{1}$, the reverse of our convention.)

Proposition 2.1. If $p_{1}, p_{2}$ and $p_{3}$ are distinct odd primes, where $p_{1}<p_{2}<p_{3}$ and where $p_{3} \equiv 1\left(\bmod p_{1}\right), p_{2} \equiv 1\left(\bmod p_{1}\right)$, but $p_{3} \not \equiv 1\left(\bmod p_{2}\right)$, then the groups of order $p_{1} p_{2} p_{3}$ are

$$
\begin{aligned}
C_{p_{3} p_{2} p_{1}} & =\left\langle x, y, z \mid x^{p_{3}}, y^{p_{2}}, z^{p_{1}}, y x y^{-1}=x, z x z^{-1}=x, z y z^{-1}=y\right\rangle, \\
C_{p_{2}} \times\left(C_{p_{3}} \rtimes C_{p_{1}}\right) & =\left\langle x, y, z \mid x^{p_{3}}, y^{p_{2}}, z^{p_{1}}, y x y^{-1}=x, z x z^{-1}=x^{v_{3}}, z y z^{-1}=y\right\rangle, \\
C_{p_{3}} \times\left(C_{p_{2}} \rtimes C_{p_{1}}\right) & =\left\langle x, y, z \mid x^{p_{3}}, y^{p_{2}}, z^{p_{1}}, y x y^{-1}=x, z x z^{-1}=x, z y z^{-1}=y^{v_{2}}\right\rangle, \\
C_{p_{3} p_{2}} \rtimes_{i} C_{p_{1}} & =\left\langle x, y, z \mid x^{p_{3}}, y^{p_{2}}, z^{p_{1}}, y x y^{-1}=x, z x z^{-1}=x^{v_{3}}, z y z^{-1}=y^{v_{2}^{v}}\right\rangle,
\end{aligned}
$$

where $i=1, \ldots, p_{1}-1$.
Our goal is to examine $R(\Gamma,[M])$ for all groups of this order, as presented above. Following the program as laid out at the end of Section 1, we shall work within the ambient symmetric group $B=S_{m p}$ where, in this case, $p=p_{3}$ and $m=p_{1} p_{2}$. Also, we will choose representative regular subgroups of $\operatorname{Norm}_{B}(\mathcal{P})$ where $\mathcal{P}$ is generated by the product of $m$ disjoint $p$-cycles. The elements of $\operatorname{Norm}_{B}(\mathcal{P})$ shall be tuples $(\hat{x}, u, \xi)$ where $\hat{x}$ is a vector in $\mathbb{F}_{p}^{m}$, and where $u \in U_{p}, \xi \in S_{m}$. We note that $\mathcal{P}$ is embedded in $\operatorname{Norm}_{B}(\mathcal{P})$ as $\langle([1,1, \ldots, 1], 1, I)\rangle$. The representation of each $\Gamma$ (as a regular subgroup of $\left.\operatorname{Norm}_{B}(\mathcal{P})\right)$ from among the $p_{1}+2$ different isomorphism classes is somewhat arbitrary but will be selected for computational convenience. Also, all will be chosen to have their Sylow $p$-subgroup be $\mathcal{P}$. The differences will lie in the representation of the complementary subgroups of order $m$, of which there are two possibilities, up to isomorphism, given that $p_{2} \equiv 1\left(\bmod p_{1}\right)$.
Lemma 2.2. If we define $\overline{v_{2}}=v_{2}^{-1}$ in $U_{p_{2}}$ and

$$
\begin{aligned}
\sigma & =\prod_{k=1}^{p_{1}}\left(k, k+p_{1}, k+2 p_{1}, \ldots, k+\left(p_{2}-1\right) p_{1}\right) \\
\tilde{\sigma} & =\prod_{k=1}^{p_{1}}\left(k, k+p_{1}, k+2 p_{1}, \ldots, k+\left(p_{2}-1\right) p_{1}\right)^{\bar{v}_{2}{ }^{k}} \\
\tau & =\prod_{i=0}^{p_{2}-1}\left(1+i p_{1}, 2+i p_{1}, \ldots, p_{1}+i p_{1}\right) \\
\delta & =\left(\prod_{i=0}^{p_{2}-1}\left(1+i p_{1}, 2+\overline{v_{2}} i p_{1}, \ldots, p_{1}+{\overline{v_{2}}}^{p_{1}-1} i p_{1}\right)\right)^{-1}
\end{aligned}
$$

then $\langle\sigma, \tau\rangle \cong C_{p_{2} p_{1}}$ and $\langle\sigma, \delta\rangle \cong C_{p_{2}} \rtimes C_{p_{1}}$, and both are regular subgroups of $S_{m}$ where $m=p_{1} p_{2}$. Moreover, $\langle\tilde{\sigma}, \tau\rangle=(\langle\sigma, \delta\rangle)^{\mathrm{opp}}=\operatorname{Cent}_{S_{m}}(\langle\sigma, \delta\rangle)$.
Proof. For the metacyclic group $C_{p_{2}} \rtimes C_{p_{1}}$, there exist generators of orders $p_{2}$ and $p_{1}$ where conjugating the order- $p_{2}$ generator by the order- $p_{1}$ generator raises the order- $p_{2}$ generator to the power $v_{2}$, where $v_{2}$ has order $p_{1}$ in $U_{p_{2}}$. This is possible given that $p_{2} \equiv 1\left(\bmod p_{1}\right)$. For $C_{p_{2} p_{1}}$, the two generators must, of course, centralize each other. If one writes $\sigma$ above as $\sigma_{1} \cdots \sigma_{p_{1}}$ then one may verify that $\tau \sigma_{i} \tau^{-1}=\sigma_{i+1}$ and that $\delta \sigma_{i} \delta^{-1}=\sigma_{i-1}^{v_{2}}$. As to regularity, one recalls that if $N$ is a semiregular subgroup of $S_{n}$ of order $n$ then $N$ is regular, which is certainly the case for the groups $\langle\sigma, \tau\rangle$ and $\langle\sigma, \delta\rangle$. The last assertion is a matter of verifying that the respective generators centralize each other, for example, that $\sigma \tilde{\sigma} \sigma^{-1}=\tilde{\sigma}$.

For the groups of order $p_{1} p_{2} p_{3}$, the generators of order $p_{2}$ centralize the order $-p_{3}$ generator, but the order- $p_{1}$ generator may or may not centralize the generators of orders $p_{2}$ and $p_{3}$ as presented in Proposition 2.1. Using these presentations, we define the following $\Gamma$ of each isomorphism class embedded in $\operatorname{Norm}_{B}(\mathcal{P})$, just as in [K, p. 2230].
Proposition 2.3. The regular subgroups of $\operatorname{Norm}_{B}(\mathcal{P})$ from each of the isomorphism classes of groups of order $m p=\left(p_{1} p_{2}\right) p_{3}$ given in Proposition 2.1 are

$$
\begin{aligned}
\Gamma=C_{p_{3} p_{2} p_{1}} & =\langle(\hat{1}, 1, I),(\hat{0}, 1, \sigma),(\hat{0}, 1, \tau)\rangle, \\
\Gamma=C_{p_{2}} \times\left(C_{p_{3}} \rtimes C_{p_{1}}\right) & =\left\langle(\hat{1}, 1, I),(\hat{0}, 1, \sigma),\left(\hat{0}, v_{3}, \tau\right)\right\rangle, \\
\Gamma=C_{p_{3}} \times\left(C_{p_{2}} \rtimes C_{p_{1}}\right) & =\langle(\hat{1}, 1, I),(\hat{0}, 1, \tilde{\sigma}),(\hat{0}, 1, \tau)\rangle, \\
\Gamma=C_{p_{3} p_{2}} \rtimes_{j} C_{p_{1}} & =\left\langle(\hat{1}, 1, I),(\hat{0}, 1, \tilde{\sigma}),\left(\hat{0}, v_{3}, \tau^{j}\right)\right\rangle,
\end{aligned}
$$

where $j=1, \ldots, p_{1}-1$.
Proof. We note that $\mathcal{P} \leq \operatorname{Norm}_{B}(\mathcal{P})$ is generated by $(\hat{1}, 1, I)$. We can prove that these groups have the asserted structure by using (2) above. First, we note that $(\hat{0}, 1, \beta)$ centralizes $(\hat{1}, 1, I)$ for any $\beta \in S_{m}$, and that $(\hat{0}, 1, \sigma)$ is centralized by $(\hat{0}, 1, \tau)$. Next we have

$$
\begin{aligned}
\left(\hat{0}, v_{3}, \tau\right)(\hat{1}, 1, I)\left(\hat{0}, v_{3}, \tau\right)^{-1} & =\left(\hat{0}+v_{3} \tau(\hat{1})-\hat{0}, v_{3} \cdot 1 \cdot v_{3}^{-1}, \tau I \tau^{-1}\right) \\
& =\left(v_{3} \tau(\hat{1}), 1, I\right)=\left(v_{3} \hat{1}, 1, I\right)=(\hat{1}, 1, I)^{v_{3}}
\end{aligned}
$$

and similarly $\left(\hat{0}, v_{3}, \tau^{j}\right)(\hat{1}, 1, I)\left(\hat{0}, v_{3}, \tau^{j}\right)^{-1}=(\hat{1}, 1, I)^{v_{3}}$. We also have

$$
(\hat{0}, 1, \tau)(\hat{0}, 1, \tilde{\sigma})(\hat{0}, 1, \tau)^{-1}=\left(\hat{0}, 1, \tau \tilde{\sigma} \tau^{-1}\right)=\left(\hat{0}, 1, \tilde{\sigma}^{v_{2}}\right)=(\hat{0}, 1, \tilde{\sigma})^{v_{2}}
$$

and $\left(\hat{0}, v_{3}, \tau^{j}\right)(\hat{0}, 1, \tilde{\sigma})\left(\hat{0}, v_{3}, \tau^{j}\right)^{-1}=(\hat{0}, 1, \tilde{\sigma})^{v_{2}^{j}}$. The proof is finished by recalling [K, Proposition 3.8], which states that if $(\hat{a}, u, \alpha)$ in $\operatorname{Norm}_{B}(\mathcal{P})$ has order coprime to $p$ then it is fixed-point-free if and only if $\alpha \in S_{m}$ is fixed-point-free. In each of
the groups listed here, the order- $p$ subgroup $\mathcal{P}$ is unique and acts freely, and the elements outside of $\mathcal{P}$ have fixed-point-free permutation coordinates. Thus each group $\Gamma$ given in the statement of the proposition is semiregular of order $m p$ and therefore regular.

Theorem 2.4. The cardinality of $R(\Gamma,[M])$ for the $p_{1}+2$ classes of groups of order $p_{1} p_{2} p_{3}$ given in Proposition 2.1 is as follows, where the rows correspond to different $\Gamma$ and the columns are the classes $[M]$ :

| $\Gamma \downarrow \quad M \rightarrow C_{p_{3} p_{2} p_{1}}$ | $C_{p_{3}} \times\left(C_{p_{2}} \rtimes C_{p_{1}}\right)$ | $C_{p_{2}} \times\left(C_{p_{3}} \rtimes C_{p_{1}}\right)$ | $C_{p_{3} p_{2} \rtimes_{i} C_{p_{1}}}$ |  |
| :---: | :---: | :---: | :---: | :---: |
| $C_{p_{3} p_{2} p_{1}}$ | 1 | $2\left(p_{1}-1\right)$ | $2\left(p_{1}-1\right)$ | $4\left(p_{1}-1\right)$ |
| $C_{p_{3}} \times\left(C_{p_{2}} \rtimes C_{p_{1}}\right)$ | $p_{2}$ | $2\left(1+p_{2}\left(p_{1}-2\right)\right)$ | $2 p_{2}\left(p_{1}-1\right)$ | $4\left(1+p_{2}\left(p_{1}-2\right)\right)$ |
| $C_{p_{2}} \times\left(C_{p_{3}} \rtimes C_{p_{1}}\right)$ | $p_{3}$ | $2 p_{3}\left(p_{1}-1\right)$ | $2\left(1+p_{3}\left(p_{1}-2\right)\right)$ | $4\left(1+p_{3}\left(p_{1}-2\right)\right.$ |
| $C_{p_{3} p_{2} \rtimes_{j} C_{p_{1}}}$ | $p_{3} p_{2}$ | $2 p_{3}\left(1+p_{2}\left(p_{1}-2\right)\right)$ | $2 p_{2}\left(1+p_{3}\left(p_{1}-2\right)\right)$ | - |

The cardinality of $R\left(\Gamma,\left[C_{p_{3} p_{2}} \rtimes_{i} C_{p_{1}}\right]\right)$ is independent of $i \in U_{p_{1}}$ for the $\Gamma$ listed on the first three rows, and for the last it depends on the relationship between $i$ and $j$ in $U_{p_{1}}$ thus:

| $i, j$ | $\left\|R\left(C_{p_{3} p_{2}} \rtimes_{j} C_{p_{1}},\left[C_{p_{3} p_{2}} \rtimes_{i} C_{p_{1}}\right]\right)\right\|$ |
| :---: | :---: |
| $j=i,-i$ | $2\left(1+p_{3}+p_{2}+\left(2 p_{1}-5\right) p_{2} p_{3}\right)$ |
| $j \neq i,-i$ | $2\left(2 p_{3}+2 p_{2}+\left(2 p_{1}-6\right) p_{2} p_{3}\right)$ |

As there are $p_{1}+2$ classes of groups of order $\left(p_{1} p_{2}\right) p_{3}$ and therefore $\left(p_{1}+2\right)^{2}$ different possible $R(\Gamma,[M])$, segmented into different classes depending on the different possibilities for $P(N)$, fully detailing the enumeration of all these would tax the patience of the reader. Moreover, given that many pairings give rise to very similar calculations, we shall instead give a sampling of the computations of $R(\Gamma,[M])$. In particular, we shall focus on the enumeration of $R\left(C_{p_{3} p_{2}} \rtimes_{j} C_{p_{1}},\left[C_{p_{3} p_{2}} \rtimes_{i} C_{p_{1}}\right]\right)$ since these can effectively be captured in one single computational framework. We will enumerate those $N$ where $P(N)=\mathcal{P}$ and double the resulting figure to account for the corresponding $N^{\mathrm{opp}}$ which arise, and therefore have the full count.

Proof. We have $\Gamma \cong C_{p_{3} p_{2}} \rtimes_{j} C_{p_{1}}$ which is embedded in $\operatorname{Norm}_{B}(\mathcal{P})$ as

$$
\left\langle(\hat{1}, 1, I),(\hat{0}, 1, \tilde{\sigma}),\left(\hat{0}, v_{3}, \tau^{j}\right)\right\rangle
$$

and we are looking at those regular $N \leq \operatorname{Norm}_{B}(\mathcal{P})$ isomorphic to $C_{p_{3} p_{2}} \rtimes_{i} C_{p_{1}}$ and normalized by this $\Gamma$, where $i, j \in U_{p_{1}}$. Moreover, as we will be focusing on those $N$ such that $P(N)=\mathcal{P}$, we have

$$
N=\left\langle(\hat{1}, 1, I),\left(\hat{a}, u_{3}^{r}, \alpha\right),\left(\hat{b}, u_{3}^{s}, \beta\right)\right\rangle,
$$

where $\left(\hat{a}, u_{3}^{r}, \alpha\right)$ has order $p_{2}$ and centralizes $(\hat{1}, 1, I)$, and $\left(\hat{b}, u_{3}^{s}, \beta\right)$ has order $p_{1}$ and conjugates $(\hat{1}, 1, I)$ to $(\hat{1}, 1, I)^{v_{3}}=\left(v_{3} \hat{1}, 1, I\right)$ and $\left(\hat{a}, u_{3}^{r}, \alpha\right)$ to $\left(\hat{a}, u_{3}^{r}, \alpha\right)^{v_{2}^{i}}$, in accordance with the presentations of the abstract groups as in Proposition 2.1. We will consider those conditions on the components of these 3-tuples which govern order and (semi)regularity and guarantee that $\Gamma$ normalizes $N$. The computations themselves will require the basic operational facts about $\operatorname{Norm}_{B}(\mathcal{P})$ as given in (1), (2), and (3) together with the fact, mentioned earlier, that if $(\hat{v}, v, \zeta)$ has order coprime to $p=p_{3}=|\mathcal{P}|$ and acts without fixed points, then $\zeta \in S_{m}\left(m=p_{2} p_{3}\right)$ must act without fixed points on the $m$ coordinates of $\hat{v}$. We shall proceed $a d$ hoc, playing off the different requirements against each other in order to limit the choices for the components.

To begin with, we have (by virtue of the isomorphism class of $N$ )

$$
\left(\hat{a}, u_{3}^{r}, \alpha\right)(\hat{1}, 1, I)\left(\hat{a}, u_{3}^{r}, \alpha\right)^{-1}=\left(u_{3}^{r} \alpha(\hat{1}), 1, \alpha I \alpha^{-1}\right)=\left(u_{3}^{r} \hat{1}, 1, I\right)=(\hat{1}, 1, I)^{u_{3}^{r}}
$$

which means that $u_{3}^{r}=1$ since the order $-p_{2}$ generator of $N$ must centralize $\mathcal{P}$, so that $\left(\hat{a}, u_{3}^{r}, \alpha\right)=(\hat{a}, 1, \alpha)$. And since $|(\hat{a}, 1, \alpha)|=p_{2}$, we have

$$
\left(\sum_{t=0}^{p_{2}-1} \alpha^{t}(\hat{a}), 1, \alpha^{p_{2}}\right)=(\hat{0}, 1, I),
$$

which, since ( $\hat{a}, 1, \alpha$ ) is fixed-point-free of order coprime to $p_{3}$, means that $\alpha$ equals $\alpha_{1} \alpha_{2} \cdots \alpha_{p_{1}}$, a product of $p_{1}$ disjoint $p_{2}$-cycles.

Also, in $N$ we have

$$
\left(\hat{b}, u_{3}^{s}, \beta\right)(\hat{1}, 1, I)\left(\hat{b}, u_{3}^{s}, \beta\right)^{-1}=\left(u_{3}^{s} \beta(\hat{1}), 1, I\right)=\left(u_{3}^{s} \hat{1}, 1, I\right)=(\hat{1}, 1, I)^{u_{3}^{s}} \text {, }
$$

which means that $u_{3}^{s}=v_{3}$, so that $\left(\hat{b}, u_{3}^{s}, \beta\right)=\left(\hat{b}, v_{3}, \beta\right)$. And since this must have order $p_{1}$, we have

$$
\left(\sum_{t=0}^{p_{1}-1} v_{3}^{t} \beta^{t}(\hat{b}), v_{3}^{p_{1}}, \beta^{p_{1}}\right)=(\hat{0}, 1, I)
$$

which, again by the fixed-point-freeness condition on this element of order coprime to $p_{3}$, means that $\beta$ is a fixed-point-free element of order $p_{1}$ in $S_{m}$.

Again in $N$, we must have, by virtue of how the order- $p_{1}$ generator must act on the order $-p_{2}$ generator, that

$$
\left(\hat{b}, v_{3}, \beta\right)(\hat{a}, 1, \alpha)\left(\hat{b}, v_{3}, \beta\right)^{-1}=\left(\hat{b}+v_{3} \beta(\hat{a})-\left(\beta \alpha \beta^{-1}\right)(\hat{b}), 1, \beta \alpha \beta^{-1}\right)
$$

where the right-hand side must equal

$$
(\hat{a}, 1, \alpha)^{v_{2}^{i}}=\left(\sum_{t=0}^{v_{2}^{i}-1} \alpha^{t}(\hat{a}), 1, \alpha^{v_{2}^{i}}\right)
$$

In particular, we have $\beta \alpha \beta^{-1}=\alpha^{\nu_{2}^{i}}$ and, more broadly, $\beta \in \operatorname{Norm}_{S_{m}}(\langle\alpha\rangle)$ where $\alpha$ is a product of $p_{1}$ disjoint $p_{2}$-cycles in $S_{m}$ where $m=p_{1} p_{2}$. The implications of this are that $\operatorname{Norm}_{S_{m}}(\langle\alpha\rangle)$ is itself another twisted wreath product just like $\operatorname{Norm}_{B}(\mathcal{P})$ and so we shall use the same sort of 3-tuple representation for understanding the relationship between $\alpha$ and $\beta$ within this smaller normalizer. We shall return to the analysis of this relationship shortly.

Looking outward, we now start imposing restrictions on the generators of $N$ imposed by $N$ being normalized by $\Gamma$. To start with, we observe that in $C_{p_{2} p_{3}} \rtimes_{i} C_{p_{1}}$ the order $-p_{2}$ and order- $p_{3}$ subgroups are characteristic since they are unique of those orders. Thus

$$
\left(\hat{0}, v_{3}, \tau^{j}\right)(\hat{a}, 1, \alpha)\left(\hat{0}, v_{3}, \tau^{j}\right)^{-1}=\left(v_{3} \tau^{j}(\hat{a}), 1, \tau^{j} \alpha \tau^{-j}\right),
$$

where the right-hand side must be an element of $\langle(\hat{a}, 1, \alpha)\rangle$; hence $\tau^{j}$ lies in $\operatorname{Norm}_{S_{m}}(\langle\alpha\rangle)$, and thus $\tau \in \operatorname{Norm}_{S_{m}}(\langle\alpha\rangle)$. We also have

$$
(\hat{0}, 1, \tilde{\sigma})(\hat{a}, 1, \alpha)(\hat{0}, 1, \tilde{\sigma})^{-1}=\left(\tilde{\sigma}(\hat{a}), 1, \tilde{\sigma} \alpha \tilde{\sigma}^{-1}\right),
$$

where the right-hand side must equal $(\hat{a}, 1, \alpha)$, since $\operatorname{Aut}\left(C_{p_{2}}\right)$ has no $p_{2}$-torsion. This implies that $\tilde{\sigma}(\hat{a})=\hat{a}$ and $\tilde{\sigma} \in \operatorname{Cent}_{S_{m}}(\alpha)$ which means that $\alpha \in \operatorname{Cent}_{S_{m}}(\tilde{\sigma})$. The latter observation implies that $\alpha \in\left\langle\sigma_{1}, \sigma_{2}, \ldots, \sigma_{p_{1}}\right\rangle$, where $\tilde{\sigma}=\sigma_{1} \sigma_{2}^{\overline{v_{2}}} \cdots \sigma_{p_{1}}^{\overline{v_{2}} p_{1}-1}$. The reason for this is that $\operatorname{Cent}_{S_{m}}(\tilde{\sigma})$ is isomorphic to the wreath product $C_{p_{2}} 2 S_{p_{1}} \cong$ $C_{p_{2}}^{p_{1}} \rtimes S_{p_{1}}$, where the base group of the wreath product, $C_{p_{2}}^{p_{1}}$, corresponds to $\left\langle\sigma_{1}, \sigma_{2}^{\overline{v_{2}}}, \ldots, \sigma_{p_{1}}^{\overline{2_{1}} p_{1}-1}\right\rangle=\left\langle\sigma_{1}, \ldots, \sigma_{p_{1}}\right\rangle$. Therefore, in $C_{p_{2}}^{p_{1}} \rtimes S_{p_{1}}$ the only $p_{2}$-torsion is in this base group since $S_{p_{1}}$ cannot have any $p_{2}$-torsion. As to $\hat{a}=\left[a_{1}, \ldots, a_{m}\right]$, the condition $\tilde{\sigma}(\hat{a})=\hat{a}$ means that, for any $k \in\{1, \ldots, m\}$,

$$
a_{k}=a_{\tilde{\sigma}(k)}=a_{\tilde{\sigma}^{2}(k)}=\cdots=a_{\tilde{\sigma}^{p_{2}-1}(k)} .
$$

But now, since $|(\hat{a}, 1, \alpha)|=p_{2}$, we have $\sum_{t=0}^{p_{2}-1} \alpha^{t}(\hat{a})=\hat{0}$, and since $\alpha$ is in $\left\langle\sigma_{1}, \sigma_{2}, \ldots, \sigma_{p_{1}}\right\rangle$, the orbit of $k \in\{1, \ldots, m\}$ under $\alpha$ is the same as under $\tilde{\sigma}$. Thus $p_{2} a_{k}$ vansishes for any $a_{k}$ in $\hat{a}$. Since $\hat{a} \in \mathbb{F}_{p_{3}}^{m}$, this means $a_{k}=0$ since $p_{2} \neq 0$ in $\mathbb{F}_{p_{3}}$. The end result is that $(\hat{a}, 1, \alpha)=(\hat{0}, 1, \alpha)$.

As we saw above, $\alpha$ belongs to $\left\langle\sigma_{1}, \sigma_{2}, \ldots, \sigma_{p_{1}}\right\rangle$, and since $\alpha$ acts freely, we must have $\alpha=\sigma_{1}^{q_{1}} \sigma_{2}^{q_{2}} \cdots \sigma_{p_{1}}^{q_{p_{1}}}$, where $q_{k} \in U_{p_{2}}$ for $k=1, \ldots, p_{1}$. Since $\tau \sigma_{i} \tau^{-1}=\sigma_{i+1}$, if $\tau \alpha \tau^{-1}=\alpha^{w}$ for some $w \in U_{p_{2}}$ then, given that $|\tau|=p_{1}$, we have $w^{p_{1}} \equiv 1\left(\bmod p_{2}\right)$ and so $w=v_{2}^{f}$ for some $f \in \mathbb{Z}_{p_{1}}$. And since

$$
\tau\left(\sigma_{1}^{q_{1}} \sigma_{2}^{q_{2}} \cdots \sigma_{p_{1}}^{q_{p_{1}}}\right) \tau^{-1}=\sigma_{1}^{q_{p_{1}}} \sigma_{2}^{q_{1}} \cdots \sigma_{p_{1}-1}^{q_{p_{1}-1}}
$$

we have $w q_{1}=q_{p_{1}}, w q_{2}=q_{1}, \ldots, w q_{p_{1}}=q_{p_{1}-1}$, which means that

$$
\hat{q}=\left[q_{1}, q_{2}, \ldots, q_{p_{1}}\right]=q_{1}\left[1, w^{p_{1}-1}, w^{p_{1}-2}, \ldots, w^{2}, w\right],
$$

where now $q_{1} w=q_{1} w^{\prime}$ if and only if $w=w^{\prime}$, and if $q_{1} \neq q_{1}^{\prime}$ then $\hat{q} \neq \hat{q}^{\prime}$, which gives, ostensibly, $\phi\left(p_{2}\right) p_{1}$ choices for $\alpha$. However, recalling that $(\hat{0}, 1, \alpha)^{k}=\left(\hat{0}, 1, \alpha^{k}\right)$, we may divide this count by $\phi\left(p_{2}\right)$ (i.e., assume $q_{1}=1$ ) to get $p_{1}$ unique choices for $\langle\alpha\rangle$ and thus $p_{1}$ unique $\langle(\hat{0}, 1, \alpha)\rangle$, one of each $w=v_{2}^{f}$ for $f \in \mathbb{Z}_{p_{1}}$.

Now that we have the enumeration of $(\hat{0}, 1, \alpha)$ given above, how many $\left(\hat{b}, v_{3}, \beta\right)$ are there? What first must be observed is that $C_{p_{2} p_{3}} \rtimes_{i} C_{p_{1}}$ has $p_{3} p_{2} \phi\left(p_{1}\right)$ elements of order $p_{1}$. In the presentation

$$
\left\langle x, y, z \mid x^{p_{3}}, y^{p_{2}}, z^{p_{1}}, y x y^{-1}=x, z x z^{-1}=x^{v_{3}}, z y z^{-1}=y^{v_{2}^{i}}\right\rangle
$$

the elements of order $p_{1}$ are of the form $x^{r} y^{s} z^{t}$ for $r \in \mathbb{Z}_{p_{3}}, s \in \mathbb{Z}_{p_{2}}, t \in U_{p_{1}}$, so the Sylow $p_{1}$-subgroups are far from characteristic. Thus, when a generator of $\Gamma$ acts on $\left(\hat{b}, v_{3}, \beta\right)$ by conjugation, the result is one of these $p_{3} p_{2} \phi\left(p_{1}\right)$ other elements of order $p_{1}$. Consider the action of conjugation by $(\hat{0}, 1, \tilde{\sigma})$ $(\hat{0}, 1, \tilde{\sigma})\left(\hat{b}, v_{3}, \beta\right)(\hat{0}, 1, \tilde{\sigma})^{-1}=\left(\tilde{\sigma}(\hat{b}), v_{3}, \tilde{\sigma} \beta \tilde{\sigma}^{-1}\right)=(\hat{1}, 1, I)^{r}(\hat{0}, 1, \alpha)^{s}\left(\hat{b}, v_{3}, \beta\right)^{t}$, which implies that $t$ must equal 1 and so

$$
(\hat{0}, 1, \tilde{\sigma})\left(\hat{b}, v_{3}, \beta\right)(\hat{0}, 1, \tilde{\sigma})^{-1}=\left(r \hat{1}+\alpha^{s}(\hat{b}), v_{3}, \alpha^{s} \beta\right)
$$

For the action of $\left(\hat{0}, v_{3}, \tau^{j}\right)$ we get

$$
\begin{aligned}
\left(\hat{0}, v_{3}, \tau^{j}\right)\left(\hat{b}, v_{3}, \beta\right)\left(\hat{0}, v_{3}, \tau^{j}\right)^{-1} & =\left(v_{3} \tau^{j}(\hat{b}), v_{3}, \tau^{j} \beta \tau^{-j}\right) \\
& =(\hat{1}, 1, I)^{r^{\prime}}(\hat{0}, 1, \alpha)^{s^{\prime}}\left(\hat{b}, v_{3}, \beta\right)^{t^{\prime}},
\end{aligned}
$$

which implies that $t^{\prime}$ must equal 1 and so

$$
(\hat{1}, 1, I)^{r^{\prime}}(\hat{0}, 1, \alpha)^{s^{\prime}}\left(\hat{b}, v_{3}, \beta\right)^{t^{\prime}}=\left(r^{\prime} \hat{1}+\alpha^{s^{\prime}}(\hat{b}), v_{3}, \alpha^{s^{\prime}} \beta\right) .
$$

This leads to four "normalization conditions" which must be satisfied:

$$
\begin{gather*}
\tilde{\sigma}(\hat{b})-\alpha^{s}(\hat{b}) \in\langle\hat{1}\rangle,  \tag{n1}\\
\tilde{\sigma} \beta \tilde{\sigma}^{-1}=\alpha^{s} \beta,  \tag{n2}\\
v_{3} \tau^{j}(\hat{b})-\alpha^{s^{\prime}}(\hat{b}) \in\langle\hat{1}\rangle,  \tag{n3}\\
\tau^{j} \beta \tau^{-j}=\alpha^{s^{\prime}} \beta . \tag{n4}
\end{gather*}
$$

We can deal with (n1) immediately by looking once more at how the generators of $N$ interact. We have

$$
\left(\hat{b}, v_{3}, \beta\right)(\hat{0}, 1, \alpha)\left(\hat{b}, v_{3}, \beta\right)^{-1}=\left(\hat{b}-\left(\beta \alpha \beta^{-1}\right)(\hat{b}), 1, \beta \alpha \beta^{-1}\right),
$$

which must equal $(\hat{0}, 1, \alpha)^{v_{2}^{i}}=\left(\hat{0}, 1, \alpha^{v_{2}^{i}}\right)$, implying that $\hat{b}=\left(\beta \alpha \beta^{-1}\right)(\hat{b})=\alpha^{v_{2}^{i}}(\hat{b})$. By a similar argument to that above, the components of $\hat{b} \in \mathbb{F}_{p_{3}}^{m}$ are also constant along the supports of $\sigma_{1}, \ldots, \sigma_{p_{1}}$, so in fact we may simply observe that $\alpha(\hat{b})=\hat{b}$.

Thus $\tilde{\sigma}(\hat{b})-\alpha^{s}(\hat{b})=\hat{b}-\hat{b}=\hat{0}$, which lies in $\langle\hat{1}\rangle$ automatically. This allows us, by the way, to rewrite (n3) as $v_{3} \tau^{j}(\hat{b})-\hat{b} \in\langle\hat{1}\rangle$. For (n2) and (n4) we must use that $\beta \in \operatorname{Norm}_{S_{m}}(\langle\alpha\rangle)$ and represent $\beta$ as an element of this normalizer. Since $\alpha=$ $\sigma_{1} \sigma_{2}^{w^{p_{1}-1}} \cdots \sigma_{p_{1}}^{w}$ (a product of $p_{1}$ disjoint $p_{2}$-cycles), given that $m=p_{1} p_{2}$ we have

$$
\operatorname{Norm}_{S_{m}}(\langle\alpha\rangle) \cong\left\langle\sigma_{1}, \sigma_{2}^{w^{p_{1}-1}}, \ldots, \sigma_{p_{1}}^{w}\right\rangle \rtimes\left(U_{p_{2}} \times S_{p_{1}}\right) \cong \mathbb{F}_{p_{2}}^{p_{1}} \rtimes\left(U_{p_{2}} \times S_{p_{1}}\right) .
$$

Thus we may write $\alpha$ as the 3-tuple ( $\hat{1}, 1, I$ ); moreover $\beta=\left(\hat{c}, v_{2}^{i}, \mu\right)$ for some $\hat{c}$ and $\mu$, since $\beta$ normalizes $\langle\alpha\rangle$ and in view of the following calculation, which uses (2) and the fact that $\mu(\hat{1})=\hat{1}$ for all $\mu$ :

$$
\left(\hat{c}, v_{2}^{i}, \mu\right)(\hat{1}, 1, I)\left(\hat{c}, v_{2}^{i}, \mu\right)^{-1}=\left(v_{2}^{i} \hat{1}, 1, I\right)=\alpha^{v_{2}^{i}} .
$$

At first glance, this permits a fairly large number of possible $\beta$ for a given $\alpha$, but the requirement that $\Gamma$ normalizes $N$ imposes quite a number of restrictions. We begin by observing that

$$
\begin{aligned}
& \alpha=(\hat{1}, 1, I), \\
& \tilde{\sigma}=\left(\left[1, w \overline{v_{2}}, w^{2}{\overline{v_{2}}}^{2}, \ldots, w^{p_{1}-1}{\overline{v_{2}}}^{p_{1}-1}\right], 1, I\right)=(\hat{d}, 1, I), \\
& \tau=\left(\hat{0}, w,\left(1,2, \ldots, p_{1}\right)\right),
\end{aligned}
$$

and, with this in mind, we see that ( n 4 ) translates into conditions on the components of these 3-tuples in $\operatorname{Norm}_{S_{m}}(\langle\alpha\rangle)$. In particular, with respect to $\beta=\left(\hat{c}, v_{2}^{i}, \mu\right)$,

$$
\begin{aligned}
\tau^{j} \beta \tau^{-j} & =\left(\hat{0}, w,\left(1, \ldots, p_{1}\right)\right)^{j}\left(\hat{c}, v_{2}^{i}, \mu\right)\left(\hat{0}, w,\left(1, \ldots, p_{1}\right)\right)^{-j} \\
& =\left(w^{j}\left(1, \ldots, p_{1}\right)^{j}(\hat{c}), v_{2}^{i},\left(1, \ldots, p_{1}\right)^{j} \mu\left(1, \ldots, p_{1}\right)^{-j}\right), \\
\alpha^{s^{\prime}} \beta & =\left(s^{\prime} \hat{1}, 1, I\right)\left(\hat{c}, v_{2}^{i}, \mu\right)=\left(s^{\prime} \hat{1}+\hat{c}, v_{2}^{i}, \mu\right) .
\end{aligned}
$$

Since $\beta=\left(\hat{c}, v_{2}^{i}, \mu\right)$ has order $p_{1}$ (and is therefore coprime to $|\alpha|$ in $\operatorname{Norm}_{S_{m}}(\langle\alpha\rangle)$ ), $\mu$ must be fixed-point-free of order $p_{1}$ in $S_{p_{1}}$ and thus a $p_{1}$-cycle. But now (n4) implies that $\left(1, \ldots, p_{1}\right)^{j} \mu\left(1, \ldots, p_{1}\right)^{-j}=\mu$, which means that $\mu=\left(1, \ldots, p_{1}\right)^{e}$ for some $e \in U_{p_{1}}$. Furthermore, (n4) also implies that $w^{j}\left(1, \ldots, p_{1}\right)^{j}(\hat{c})-\hat{c} \in\langle\hat{1}\rangle$, a condition which we shall get back to shortly. Since

$$
\begin{aligned}
\tilde{\sigma} \beta \tilde{\sigma} & =(\hat{d}, 1, I)\left(\hat{c}, v_{2}^{i}, \mu\right)(\hat{d}, 1, I)^{-1}=\left(\hat{d}+\hat{c}-v_{2}^{i} \mu(\hat{d}), v_{2}^{i}, \mu\right), \\
\alpha^{s} \beta & =\left(s \hat{1}+\hat{c}, v_{2}^{i}, \mu\right),
\end{aligned}
$$

condition (n2) implies that $\hat{d}-v_{2}^{i} \mu(\hat{d}) \in\langle\hat{1}\rangle$. (All instances of $\hat{1}$ here refer to the vector $[1, \ldots, 1] \in \mathbb{F}_{p_{2}}^{p_{1}}$ which is the base group for the twisted wreath product $\left.\operatorname{Norm}_{S_{m}}(\langle\alpha\rangle).\right)$

Recalling that $\hat{d}=\left[1, w \bar{v}_{2}, w^{2} \bar{v}_{2}^{2}, \ldots, w^{p_{1}-1}{\overline{v_{2}}}^{p_{1}-1}\right]$ and that $\mu=\left(1, \ldots, p_{1}\right)^{e}$, the condition $\hat{d}-v_{2}^{i} \mu(\hat{d}) \in\langle\hat{1}\rangle$ can be analyzed by looking at the components and observing that, in $\langle\hat{1}\rangle$, all the components of a given vector are equal. That is,

$$
\begin{aligned}
\hat{d} & =\left[1, w \bar{v}_{2}, w^{2}{\overline{v_{2}}}^{2}, \ldots, w^{p_{1}-1}{\overline{v_{2}}}^{p_{1}-1}\right] \\
v_{2}^{i} \mu(\hat{d}) & =\left[\left(w \overline{v_{2}}\right)^{p_{1}-e} v_{2}^{i},\left(w \overline{v_{2}}\right)^{p_{1}-e+1} v_{2}^{i}, \ldots,\left(w \overline{v_{2}}\right)^{p_{1}-e+\left(p_{1}-1\right)} v_{2}^{i}\right] .
\end{aligned}
$$

In particular, the difference of the first components of $\hat{d}$ and $v_{2}^{i} \mu(\hat{d})$ equals the difference of their second components, and so

$$
\left[1-w^{p_{1}-e}{\overline{v_{2}}}^{p_{1}-e-i}\right]=w \bar{v}_{2}\left[1-w^{p_{1}-e}{\overline{v_{2}}}^{p_{1}-e-i}\right]
$$

in $\mathbb{F}_{p_{2}}$. If we let $x=1-w^{p_{1}-e} \overline{v_{2}}{ }^{p_{1}-e-i}$ then the above implies that $x=w \overline{v_{2}} x$, so either $x=0$, or $w \overline{v_{2}}=1$ regardless of $x$. If $w \overline{v_{2}}=1$, then, since (as determined above) $w=v_{2}^{f}$ for some $f \in \mathbb{Z}_{p_{1}}$, it must be that $f=1$. Otherwise, if $x=0$ then $w^{p_{1}-e}=v_{2}^{p_{1}-e-i}$ and, since (as determined above) $w=v_{2}^{f}$, this means that $f e \equiv e+i\left(\bmod p_{1}\right)$, which, by the way, is impossible if $f=1$ since $i \neq 0$.

Thus, if $f=1$ then there are no restrictions on $e \in U_{p_{1}}$, and if $f \neq 1$ then $f e=e+i$ which implies that $e=i(f-1)^{-1}\left(\bmod p_{1}\right)$.

Now if we go back to (n4), we have the condition $w^{j}\left(1, \ldots, p_{1}\right)^{j}(\hat{c})-\hat{c} \in\langle\hat{1}\rangle$ which means that

$$
\left[c_{1}, \ldots, c_{p_{1}}\right]=[l, l, \ldots, l]+w^{j}\left[c_{-j+1}, c_{-j+2}, \ldots, c_{-j+p_{1}}\right]
$$

for some $l \in \mathbb{F}_{p_{2}}$. Looking at the coordinates of $\hat{c}$ we get

$$
\begin{aligned}
c_{1+j} & =l+w^{j} c_{1} \\
c_{1+2 j} & =l+w^{j} c_{1+j}=l\left(1+w^{j}\right)+w^{2 j} c_{1} \\
& \vdots \\
c_{1+k j} & =l\left(1+w^{j}+\left(w^{j}\right)^{2}+\cdots+\left(w^{j}\right)^{k-1}\right)+w^{k j} c_{1} \\
& = \begin{cases}l k+c_{1} & \text { if } w=1(\text { i.e., } f=0) \\
l\left(\frac{1-\left(w^{j}\right)^{k}}{1-w^{j}}\right)+w^{k j} c_{1} & \text { if } w \neq 1 \text { (i.e., } f \neq 0)\end{cases}
\end{aligned}
$$

This gives a partial parametrization of the possible $\hat{c}$, but we must also include the conditions imposed by the fact that $|\beta|=\left|\left(\hat{c}, v_{2}^{i}, \mu\right)\right|=p_{1}$, that is,

$$
\begin{equation*}
\sum_{t=0}^{p_{1}-1} v_{2}^{i t} \mu^{t}(\hat{c})=\hat{0} \tag{*}
\end{equation*}
$$

which, since $\mu^{t}=\left(1, \ldots, p_{1}\right)^{e t}$, means that $\sum_{t=0}^{p_{1}-1} v_{2}^{i t}\left(1, \ldots, p_{1}\right)^{e t}(\hat{c})=\hat{0}$. So if we let $\hat{c}=\left[c_{1}, \ldots, c_{p_{1}}\right]$ then $\left(1, \ldots, p_{1}\right)^{e t}(\hat{c})=\left[c_{-e t+1}, c_{-e t+2}, \ldots, c_{-e t+p_{1}}\right]$, which translates into the (single) condition

$$
\begin{equation*}
\sum_{t=0}^{p_{1}-1} v_{2}^{i t} c_{-e t+1}=0 \tag{**}
\end{equation*}
$$

since the vector equation $(*)$ consists of a system of equations, all of which are
equivalent to $(* *)$. To utilize this information, together with the above parametrization of $\hat{c}$ (in terms of $c_{1}$ and ' $l$ ' above), we first observe that $1+k j=-e t+1$ implies $k=j^{-1}(-e t)$ and so

$$
c_{-e t+1}= \begin{cases}l\left(j^{-1}(-e t)\right)+c_{1} & \text { if } w=1 \text { i.e., } f=0 \\ l\left(\frac{1-w^{-e t}}{1-w^{j}}\right)+w^{-e t} c_{1} & \text { if } w \neq 1 \text { i.e., } f \neq 0 .\end{cases}
$$

For the case $w=1(f=0)$, we have $e=i(-1)^{-1}=-i$ and so

$$
\begin{aligned}
\sum_{t=0}^{p_{1}-1} v_{2}^{i t} c_{-e t+1} & =\sum_{t=0}^{p_{1}-1} v_{2}^{i t}\left(l\left(j^{-1}(-e t)\right)+c_{1}\right)=\sum_{t=0}^{p_{1}-1} v_{2}^{i t}\left(l\left(j^{-1}(i t)\right)+c_{1}\right) \\
& =\sum_{t=0}^{p_{1}-1} v_{2}^{i t}\left(l j^{-1} i t+c_{1}\right)=c_{1} \sum_{t=0}^{p_{1}-1} v_{2}^{i t}+l \sum_{t=0}^{p_{1}-1} j^{-1} i t v_{2}^{i t} \\
& =l j^{-1} i \sum_{t=0}^{p_{1}-1} t v_{2}^{i t}=l j^{-1} i \frac{p_{1}}{v_{2}^{i t}-1} .
\end{aligned}
$$

The last two lines of the above calculation are justified as follows. Since $v_{2}^{p_{1}}=1$ in $\mathbb{F}_{p_{2}}$, we have

$$
\sum_{t=0}^{p_{1}-1} v_{2}^{i t}=0, \quad \sum_{t=0}^{p_{1}-1} t x^{t}=x\left(\frac{p_{1} x^{p_{1}-1}}{x-1}-\frac{x^{p_{1}}-1}{(x-1)^{2}}\right)
$$

and substituting in $x=v_{2}^{i}$ we get

$$
\sum_{t=0}^{p_{1}-1} t v_{2}^{i t}=\frac{p_{1} v_{2}^{i p_{1}}}{v_{2}^{i}-1}=\frac{p_{1}}{v_{2}^{i}-1}
$$

This being the case, $\sum_{t=0}^{p_{1}-1} v_{2}^{i t} c_{-e t+1}=0$ if and only if $l=0$, which means that $c_{1+k j}=c_{1}$ for all $k$, and therefore $\hat{c} \in\langle\hat{1}\rangle$.

For the case where $w \neq 1$ (i.e., $f \neq 0$ ), we have

$$
\begin{aligned}
\sum_{t=0}^{p_{1}-1} v_{2}^{i t} c_{-e t+1} & =\sum_{t=0}^{p_{1}-1} v_{2}^{i t}\left(l\left(\frac{1-\left(w^{j}\right)^{j-1}(-e t)}{1-w^{j}}\right)+w^{-e t} c_{1}\right) \\
& =\frac{l}{1-w^{j}} \sum_{t=0}^{p_{1}-1} v_{2}^{i t}\left(1-w^{-e t}\right)+\sum_{t=0}^{p_{1}-1} v_{2}^{i t} w^{-e t} c_{1} \\
& =\frac{l}{1-w^{j}} \sum_{t=0}^{p_{1}-1} v_{2}^{i t}-\frac{l}{1-w^{j}} \sum_{t=0}^{p_{1}-1} v_{2}^{i t} w^{-e t}+c_{1} \sum_{t=0}^{p_{1}-1} v_{2}^{i t} w^{-e t}
\end{aligned}
$$

$$
\begin{aligned}
& =\frac{-l}{1-w^{j}} \sum_{t=0}^{p_{1}-1} v_{2}^{i t+f(-e t)}+c_{1} \sum_{t=0}^{p_{1}-1} v_{2}^{i t+f(-e t)} \\
& =\frac{-l}{1-w^{j}} \sum_{t=0}^{p_{1}-1} v_{2}^{i t-f e t}+c_{1} \sum_{t=0}^{p_{1}-1} v_{2}^{i t-f e t} \\
& =\left(c_{1}-\frac{l}{1-w^{j}}\right) \sum_{t=0}^{p_{1}-1} v_{2}^{(i-f e) t}
\end{aligned}
$$

If $i \neq f e$ then the last sum above equals 0 for all $c_{1}, l \in \mathbb{F}_{p_{2}}$, yielding $p_{2}^{2}$ choices. If $i=f e$ then the last sum is 0 only when $c_{1}=l /\left(1-w^{j}\right)$, which means only $p_{2}$ choices.

The requirements of condition (n3), that $v_{3} \tau^{j}(\hat{b})-\hat{b} \in\langle\hat{1}\rangle$, demand that one use the fact seen earlier, namely that $\tilde{\sigma}(\hat{b})=\hat{b}$ and equivalently $\alpha(\hat{b})=\hat{b}$. We also must factor in order considerations, just as in the above enumeration of $\hat{c}$, namely that $\sum_{t=0}^{p_{1}-1} v_{3}^{t} \beta^{t}(\hat{b})=\hat{0}$. Since the components of $\hat{b}$ (a vector in $\mathbb{F}_{p_{3}}^{p_{1} p_{2}}$ ) are equal on the supports of the cycles that make up $\tilde{\sigma}$ (a product of $p_{1}$ disjoint $p_{2}$-cycles) and since $\tau \sigma_{i} \tau^{-1}=\sigma_{i+1}$, by identifying together these identical components, we can proceed, for the moment, as if $\hat{b}$ were a vector in $\mathbb{F}_{p_{3}}^{p_{1}}$. With this identification, $\tau$ acts on this $\hat{b}$ as $\left(1, \ldots, p_{1}\right)$ and therefore $\tau^{j}$ acts like $\left(1, \ldots, p_{1}\right)^{j}$. Similarly, since $\beta=\left(\hat{c}, v_{2}^{i},\left(1, \ldots, p_{1}\right)^{e}\right)$, it acts on $\hat{b}$ as $\left(1, \ldots, p_{1}\right)^{e}$. Consequently, if we set $\hat{b}=\left[b_{1}, b_{2}, \ldots, b_{p_{1}}\right]$ then (n3) implies that $\left[b_{1}, \ldots, b_{p_{1}}\right]=$ $[l, l, \ldots, l]+v_{3}\left[b_{-e+1}, b_{-e+2}, \ldots, b_{-e}\right]$ for some $l \in \mathbb{F}_{p_{3}}$ so that, in a similar fashion to the computation of $\hat{c}$ a few pages back, we have

$$
b_{1+k j}=l\left(1+v_{3}+v_{3}^{2}+\cdots+v_{3}^{k-1}\right)+v_{3}^{k} b_{1}=l\left(\frac{1-v_{3}^{k}}{1-v_{3}}\right)+v_{3}^{k} b_{1}
$$

Since $1+k j=1-$ et implies $k=-j^{-1}$ et, we get

$$
b_{1-e t}=l\left(\frac{1-v_{3}^{-j^{-1} e t}}{1-v_{3}}\right)+v_{3}^{-j^{-1} e t} b_{1}
$$

for the parametrization of $\hat{b}$. The question is: how many 'degrees of freedom' do we have since, ostensibly, we can choose $l, b_{1} \in \mathbb{F}_{p_{3}}$ ? The order requirement becomes $\sum_{t=0}^{p_{1}-1} v_{3}^{t}\left(1, \ldots, p_{1}\right)^{e t}(\hat{b})=\hat{0}$ which reduces to $\sum_{t=0}^{p_{1}-1} v_{3}^{t} b_{1-e t}=0$. In a similar fashion to the calculations above, we get that $\left|\left(\hat{b}, v_{3}, \beta\right)\right|=p_{1}$ if and only if

$$
\left(b_{1}-\frac{l}{1-v_{3}}\right) \sum_{t=0}^{p_{1}-1} v_{3}^{t\left(1-j^{-1} e\right)}=0
$$

which comes down to two possibilities given that $v_{3}^{p_{1}} \equiv 1\left(\bmod p_{3}\right)$. If $e \neq j$ then we may choose $b_{1}$ and $l$ in $\mathbb{F}_{p_{3}}$ arbitrarily (i.e., $p_{3}^{2}$ choices); otherwise, one must choose $b_{1}$ and $l$ such that $b_{1}=l /\left(1-v_{3}\right)$, which yields $p_{3}$ choices.

The enumeration of the possible $\left(\hat{b}, v_{3}, \beta\right)$ comes down to the interaction between the parameters $f, e, i$, and $j$ as determined by order conditions on $\left(\hat{b}, v_{3}, \beta\right)$, where $\beta=\left(\hat{c}, v_{2},\left(1, \ldots, p_{1}\right)^{e}\right)$, and the normalization conditions (n1)-(n4). The parameters $i$ and $j$ are chosen at the outset, but the core parameter is $f \in \mathbb{Z}_{p_{1}}$ which determines the possible $\alpha$. Subsequently, $e$ is determined by $\beta$ since $\left(\hat{b}, v_{3}, \beta\right)$ normalizes $\langle(\hat{0}, 1, \alpha)\rangle$. The different possibilities are summarized as follows:

$$
\begin{aligned}
& f=0 \text { implies } e=i(f-1)^{-1}=-i \\
& f=1 \text { allows } e=1, \ldots, p_{1}-1 \\
& f=2, \ldots, p_{1}-1 \text { implies } e=i(f-1)^{-1} \\
& f=0 \text { implies } p_{2} \text { choices for } \hat{c} \\
& f>0 \text { implies }\left\{\begin{array}{l}
p_{2} \text { choices for } \hat{c} \text { when } i=f e \\
p_{2}^{2} \text { choices for } \hat{c} \text { when } i \neq f e \\
j=e \text { implies } p_{3} \text { choices of } \hat{b} \\
j \neq e \text { implies } p_{3}^{2} \text { choices of } \hat{b}
\end{array}\right.
\end{aligned}
$$

If we denote by $s_{0}, s_{1}, s_{>1}$ the number of $\left(\hat{b}, v_{3}, \beta\right)$ for the different choices of $f$, then we want to know $s_{0}+s_{1}+s_{>1}$.

For $f=0$ we have $e=-i$ and so there are $p_{2}$ different $\beta=\left(\hat{c}, v_{3},\left(1, \ldots, p_{1}\right)^{e}\right)$. If $j=-i=e$ there are $p_{3}$ choices for $\hat{b}$, and if $j \neq-i$ there are $p_{3}^{2}$. Hence

$$
s_{0}= \begin{cases}p_{2} p_{3} & j=-i, \\ p_{2} p_{3}^{2} & j \neq-i\end{cases}
$$

For $f=1$ we have $e=1, \ldots, p_{1}-1$ and $f e=f$, so $f e=i$ for exactly one $e$ and $j=e$ also exactly once. Depending on whether $i=j$ or not, for potentially the same $e$, this results in different possibilities for the number of $\hat{c}$ and $\hat{b}$. We have

$$
s_{1}= \begin{cases}p_{2} p_{3}+\left(p_{1}-2\right) p_{2}^{2} p_{3}^{2} & j=i, \\ p_{2} p_{3}^{2}+p_{2}^{2} p_{3}+\left(p_{1}-3\right) p_{2}^{2} p_{3}^{2} & j \neq i\end{cases}
$$

For $f>1$ we have $e=i(f-1)^{-1}$ and so fe equals if $(f-1)^{-1}$ which is never equal to $i$; thus there are $p_{2}^{2}$ choices for $\hat{c}$. If $j=i(f-1)^{-1}=e$ then $f-1=i j^{-1}$ which is impossible if $f-1=-1$, that is, $j=-i$; thus there are $p_{3}^{2}$ different $\hat{b}$ for each $f>1$. We have then the count for $f>1$ :

$$
s_{>1}= \begin{cases}\left(p_{1}-2\right) p_{2}^{2} p_{3}^{2} & j=-i \\ \left(p_{1}-3\right) p_{2}^{2} p_{3}^{2}+p_{2}^{2} p_{3} & j \neq-i\end{cases}
$$

Now, for $i, j \in U_{p_{1}}$, we have that $j=i$ implies $j \neq-i$ since $p_{1}>2$ and, similarly, if $j=-i$ then $j \neq i$. So we have

$$
s_{0}+s_{1}+s_{>1}= \begin{cases}p_{2} p_{3}+p_{2} p_{3}^{2}+p_{2}^{2} p_{3}+\left(2 p_{1}-5\right) p_{2}^{2} p_{3}^{2} & j=i,-i \\ 2 p_{2} p_{3}^{2}+2 p_{2}^{2} p_{3}+\left(2 p_{1}-6\right) p_{2}^{2} p_{3}^{2} & j \neq i,-i\end{cases}
$$

What $s_{0}+s_{1}+s_{>1}$ represents is the number of those

$$
\left\{(\hat{0}, 1, \alpha),\left(\hat{b}, v_{2}, \beta\right)\right\}
$$

which generate $Q(N)$, where $P(N)=\mathcal{P}$ and where, for distinct $f$, the resulting $\langle\alpha\rangle$ and thus $\langle(\hat{0}, 1, \alpha)\rangle$ are distinct. Thus, to remove duplicate $Q(N)$ s (arising from $\left(\hat{b}, v_{3}, \beta\right)$ which generate the same $Q(N)$ with $\left.(\hat{0}, 1, \alpha)\right)$ we must divide by $p_{2} p_{3}$. The reason for this is that, as we mentioned above, in the abstract groups $C_{p_{2} p_{3}} \rtimes_{i} C_{p_{1}}$, if one multiplies a given element of order $p_{1}$ by an element of order $p_{2}$ or $p_{3}$ (or both) one gets another element of order $p_{1}$.

We have now completely enumerated those $N \in R\left(C_{p_{2} p_{3}} \rtimes_{j} C_{p_{1}},\left[C_{p_{2} p_{3}} \rtimes_{i} C_{p_{1}}\right]\right)$ where $P(N)=\mathcal{P}$. Since the order- $p_{3}$ subgroup is not a direct factor, we now double this figure since the groups in $R\left(C_{p_{2} p_{3}} \rtimes_{j} C_{p_{1}},\left[C_{p_{2} p_{3}} \rtimes_{i} C_{p_{1}}\right]\right)$ are evenly distributed between those classes where $P(N)=\mathcal{P}$ versus those for which $P(N) \neq \mathcal{P}$. The count given in the statement of the theorem for $\left|R\left(C_{p_{2} p_{3}} \rtimes_{j} C_{p_{1}},\left[C_{p_{2} p_{3}} \rtimes_{i} C_{p_{1}}\right]\right)\right|$ is now verified.

## 3. Square-free groups where $\boldsymbol{p}<\boldsymbol{m}$

There are many prime triples $\left(p_{1}, p_{2}, p_{3}\right)$, where $p_{3} \equiv 1\left(\bmod p_{1}\right), p_{2} \equiv 1\left(\bmod p_{1}\right)$, and $p_{3} \equiv \equiv 1\left(\bmod p_{2}\right)($ which give rise to groups of the type studied in Theorem 2.4) but where $p=p_{3}<p_{1} p_{2}=m$. Indeed, if one takes prime triples from $\{2, \ldots, 113\}$ then, of these, 474 have the property implying that groups of order $p_{1} p_{2} p_{3}$ are in the category studied in Theorem 2.4, and, of these, 246 have the property $p<m$. If we look beyond to groups of order $p_{1} p_{2} p_{3} p_{4}$, where $p_{1}<p_{2}<p_{3}<p_{4}$, which are also explored in [Alonso 1976], then the analog of Theorem 2.4 is the case where $\left\{p_{4}, p_{3}, p_{2}\right\}$ are all equivalent to $1\left(\bmod p_{1}\right)$ but none of $\left\{p_{4}, p_{3}, p_{2}\right\}$ are equivalent to 1 mod each other. In this case, the number of groups of order $p_{1} p_{2} p_{3} p_{4}$ is $p_{1}^{2}+p_{1}+2$. If one looks at the 4 -tuples of distinct primes chosen in $\{2, \ldots, 113\}$ then, of these, 3173 satisfy the congruence conditions of this class, and if $m=p_{1} p_{2} p_{3}$ and $p=p_{4}$ then, of these, 3151 have the property that $p<m$.

## Acknowledgement

The author wishes to thank the referee for the careful reading given to this paper and for the helpful comments and suggestions.

## References

[Alonso 1976] J. Alonso, "Groups of square-free order, an algorithm", Math. Comp. 30:135 (1976), 632-637. MR 58 \#22295 Zbl 0335.20002
[Byott 2004] N. P. Byott, "Hopf-Galois structures on Galois field extensions of degree pq", J. Pure Appl. Algebra 188:1-3 (2004), 45-57. MR 2004j:16041 Zbl 1047.16022
[Chase and Sweedler 1969] S. U. Chase and M. E. Sweedler, Hopf algebras and Galois theory, Lecture Notes in Mathematics 97, Springer, Berlin-New York, 1969. MR 41 \#5348 Zbl 0197.01403
[Childs 2003] L. N. Childs, "On Hopf Galois structures and complete groups", New York J. Math. 9 (2003), 99-115. MR 2004k:16097 Zbl 1038.12003
[Greither and Pareigis 1987] C. Greither and B. Pareigis, "Hopf Galois theory for separable field extensions", J. Algebra 106:1 (1987), 239-258. MR 88i:12006 Zbl 0615.12026
[Hölder 1895] O. Hölder, "Die Gruppen mit quadratfreier Ordnungzahl", Nachr. Königl. Gesell. Wissenschaft. Göttingen Math.-Phys. Kl. (1895), 211-229.
[Kohl 2013] T. Kohl, "Regular permutation groups of order $m p$ and Hopf Galois structures", Algebra Number Theory 7:9 (2013), 2203-2240. MR 3152012 Zbl 1286.12002
[Pakianathan and Shankar 2000] J. Pakianathan and K. Shankar, "Nilpotent numbers", Amer. Math. Monthly 107:7 (2000), 631-634. MR 2001i:20051 Zbl 0986.20026

Communicated by Susan Montgomery
Received 2014-08-05 Revised 2015-10-01 Accepted 2015-11-27
tkohl@math.bu.edu Department of Mathematics and Statistics, Boston University, 111 Cummington Mall, Boston, MA 02215, United States

# Algebra \& Number Theory 

msp.org/ant

## EDITORS

Managing Editor
Bjorn Poonen
Massachusetts Institute of Technology
Cambridge, USA

Editorial Board Chair
David Eisenbud
University of California
Berkeley, USA

## Board of Editors

| Georgia Benkart | University of Wisconsin, Madison, USA | Susan Montgomery | University of Southern California, USA |
| ---: | :--- | ---: | :--- |
| Dave Benson | University of Aberdeen, Scotland | Shigefumi Mori | RIMS, Kyoto University, Japan |
| Richard E. Borcherds | University of California, Berkeley, USA | Raman Parimala | Emory University, USA |
| John H. Coates | University of Cambridge, UK | Jonathan Pila | University of Oxford, UK |
| J-L. Colliot-Thélène | CNRS, Université Paris-Sud, France | Anand Pillay | University of Notre Dame, USA |
| Brian D. Conrad | Stanford University, USA | Victor Reiner | University of Minnesota, USA |
| Hélène Esnault | Freie Universität Berlin, Germany | Peter Sarnak | Princeton University, USA |
| Hubert Flenner | Ruhr-Universität, Germany | Joseph H. Silverman | Brown University, USA |
| Sergey Fomin | University of Michigan, USA | Michael Singer | North Carolina State University, USA |
| Edward Frenkel | University of California, Berkeley, USA | Vasudevan Srinivas | Tata Inst. of Fund. Research, India |
| Andrew Granville | Université de Montréal, Canada | J. Toby Stafford | University of Michigan, USA |
| Joseph Gubeladze | San Francisco State University, USA | Ravi Vakil | Stanford University, USA |
| Roger Heath-Brown | Oxford University, UK | Michel van den Bergh | Hasselt University, Belgium |
| Craig Huneke | University of Virginia, USA | Marie-France Vignéras | Université Paris VII, France |
| Kiran S. Kedlaya | Univ. of California, San Diego, USA | Kei-Ichi Watanabe | Nihon University, Japan |
| János Kollár | Princeton University, USA | Efim Zelmanov | University of California, San Diego, USA |
| Yuri Manin | Northwestern University, USA | Shou-Wu Zhang | Princeton University, USA |
| Philippe Michel | École Polytechnique Fédérale de Lausanne |  |  |

## PRODUCTION

production@msp.org
Silvio Levy, Scientific Editor

See inside back cover or msp.org/ant for submission instructions.
The subscription price for 2016 is US $\$ 290 /$ year for the electronic version, and $\$ 485 /$ year ( $+\$ 55$, if shipping outside the US) for print and electronic. Subscriptions, requests for back issues and changes of subscribers address should be sent to MSP.

Algebra \& Number Theory (ISSN 1944-7833 electronic, 1937-0652 printed) at Mathematical Sciences Publishers, 798 Evans Hall \#3840, c/o University of California, Berkeley, CA 94720-3840 is published continuously online. Periodical rate postage paid at Berkeley, CA 94704, and additional mailing offices.

ANT peer review and production are managed by EditFLOw ${ }^{\circledR}$ from MSP.

## PUBLISHED BY <br> mathematical sciences publishers <br> nonprofit scientific publishing

http://msp.org/
© 2016 Mathematical Sciences Publishers

## Algebra \& Number Theory

Volume 10 No. 12016
Stable sets of primes in number fields ..... 1ALEXANDER IVANOV
Hopf-Galois structures arising from groups with unique subgroup of order $p$ ..... 37
Timothy Kohl
On tensor factorizations of Hopf algebras ..... 61Marc Keilberg and Peter Schauenburg
Extension theorems for reductive group schemes ..... 89Adrian Vasiu
Actions of some pointed Hopf algebras on path algebras of quivers ..... 117Ryan Kinser and Chelsea Walton
On the image of the Galois representation associated to a non-CM Hida family ..... 155
Jaclyn Lang
Linear relations in families of powers of elliptic curves ..... 195
Fabrizio Barroero and Laura Capuano


[^0]:    MSC2010: primary 20B35; secondary 20D20, 20D45, 16 T 05.
    Keywords: Hopf-Galois extension, regular subgroup.

